

KULLBACK-LEIBLER INFORMATION FUNCTION AND THE SEQUENTIAL
SELECTION OF EXPERIMENTS TO DISCRIMINATE AMONG
SEVERAL LINEAR MODELS

Abstract

by

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Assume that a finite set of potential linear models relating several controlled variables to an observed variable is postulated and that exactly one of these models is the true model. The problem is to sequentially design most informative experiments so that the correct model can be determined with as little experimentation as possible. We assume that the error variance of the process is known. In addition, we assume the statistician possesses prior information which can be expressed as the prior probability that each of the proposed models is indeed the correct model and prior multivariate normal distributions on the parameters of each of the postulated model equations. After each stage of sampling, the prior distributions and the observed data values are used to compute posterior probabilities of the models being the true one and posterior distributions on the parameters of the models. Then sampling is terminated if either a prespecified number of observations has been taken or if any of the posterior probabilities of the models exceeds a prespecified minimum stopping probability. Upon termination of sampling, the model with the largest posterior

probability is chosen to be the correct model. If sampling is not to be terminated, the next experiment chosen is that one in the set of allowable values of the controlled variables which maximizes the expected Kullback-Leibler information function based upon the current posterior probabilities and distributions.

An analytical study of this procedure is too complex and difficult to adequately achieve. Hence a number of Monte-Carlo simulation experiments were performed to obtain information about the performance of this adaptive design procedure. Two basic types of Monte-Carlo experiments were performed. In the first, one of the models was chosen to be used to generate the random observations using known fixed values for the parameters. Then a large number of observations were taken using the Kullback-Leibler information functions as a criterion to choose the sequence of experiments. It was found the posterior probability of the chosen model relatively rapidly approaches the value of 1.0 and then fluctuates near 1.0. The posterior mean of the parameters of the correct model also rapidly approaches the known fixed values used to generate the observations. In the second type of experiment, one of the models was chosen to be used to generate the random observations. Then for various combinations of the maximum number of observations, stopping probability, prior distributions of the parameters, and error variance of the process, a large number of repetitions of the sequential design procedure were executed. Then a probability of correct selection and average sample number were calculated based upon the

number of times the procedure chose the correct model and the number of observations taken until termination. In general, it was found that as long as the prior mean of the correct model is not too distant from the true value with respect to the means of the other models the probability of correct selection is respectably high.

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CHAPTER 1 - INTRODUCTION AND LITERATURE SURVEY

The general linear model has become one of the most useful statistical tools available to the modern scientific experimenter. There have been many books and papers written about techniques for choosing the appropriate or "best" linear model to fit to a set of data already collected. In general, these have been methods of hypothesis testing to determine which of a set of specified terms in a model equation may be dropped from the model. Much work has also been done with regard to the problem of designing best or optimal experiments to estimate the parameters of specified model equations.

In this dissertation we study a sequential adaptive experimental design procedure for a related problem. Assume that a finite set of potential linear models relating certain controlled variables to an observed variable is postulated and that exactly one of these models is correct. The problem is to sequentially design most informative experiments so that the correct model equation can be determined with as little experimentation as possible. We also assume that the error variance of the process is known. In addition, we assume that the statistician possesses prior information which can be expressed by the prior probability that each of the proposed models is indeed the correct model and prior multivariate normal distributions on the parameters of the various models. We then de-

rive an adaptive procedure for designing the successive experiments using the Kullback-Leibler information function to maximize the anticipated information for discriminating among the models. That is, after each stage of sampling, the prior distributions and the observed values are used to compute posterior probabilities of the postulated models being correct and posterior distributions on the parameters of the models. Then if sampling is not to be terminated, the next experiment chosen is that which maximizes the expected Kullback-Leibler information based on the current posterior probabilities and distributions. Sampling is terminated whenever either a prespecified number of observations is finally taken or whenever any of the posterior probabilities of the models exceeds a prespecified value. Upon termination of sampling, the model with the largest posterior probability is chosen to be the correct model.

An analytical study of this procedure is too complex and difficult to adequately achieve. Hence a number of Monte-Carlo simulation experiments were performed to obtain information about the performance of this adaptive design procedure. Two basic types of Monte-Carlo experiments were performed. In the first, one of the models was chosen to be used to generate the random observations using known fixed values for the parameters. Then a large number of observations were taken using the Kullback-Leibler information as a criterion to choose the sequence of experiments. It was found the posterior probability of the chosen model relatively rapidly approaches the value of 1.0 and then fluctuates near 1.0. The

posterior mean of the parameters of the correct model also rapidly approach the known fixed values used to generate the observations. In the second type of experiment, one of the models was chosen to be used to generate the random observations. Then for various combinations of the maximum number of observations, stopping probability, prior distributions of the parameters, and error variance of the process, a large number of repetitions of the sequential design procedure were executed. Then a probability of correct selection and average sample number were calculated based upon the number of times the procedure chose the correct model and the number of observations taken until termination. In general, it was found that as long as the prior mean of the correct model is not too distant from the true value with respect to the means of the other models the probability of correct selection is respectably high.

We now briefly indicate the general organization of the dissertation. In Chapter 2 the notation used is described and the structure of the linear models is derived. Chapter 3 then develops the distribution theory which will be basic to the remainder of the dissertation. In particular, the posterior probabilities of the models, the posterior distributions of the parameters, and the Markovian nature of the sampling process are developed. Some large sample results are then derived for the situation where the sequence of experiments is specified in advance of experimenting. These results do not thus formally apply to the adaptive design procedure.

We find, however, that they do appear to be true to a surprising extent and provide some help in explaining and interpreting the Monte-Carlo results.

In Chapter 4, the Kullback-Leibler information concept is introduced and the derivation of the anticipated information as a function of the current posterior probabilities of the models and the current posterior distributions of the parameters is presented. This anticipated information is the criterion function used to define the most informative experiment. Its use is discussed both from the point of view of its relation to the expected decrease in entropy and the point of view that it results in a very simple function measuring the amount by which the expected value of the observed variable under each model is separated.

The sequential experiment selection, stopping, and model selection rules are presented in Chapter 5.

In Chapter 6, the Monte-Carlo simulation experiments are described and the results presented and discussed. Chapter 7 presents an example of application. Several appendixes are also included. Of most importance is appendix A which presents the computer program used to perform the simulation experiments.

We now turn to a discussion of works by earlier authors who have considered similar problems.

Lindley (1956) was one of the first to consider the general idea of applying information concepts to the problems of statistical inference. He modified the concept of entropy and developed a num-

ber of interesting general results on the amount of information in an experiment about the parameters of the distribution of a random variable.

Stone (1959) was one of the first to consider information concepts as applied to designing and comparing regression experiments. He used a Bayesian framework, but the problem he considers is that of parameter estimation rather than that of model selection.

Another early and more relevant paper is that of Chernoff (1959) who applied the Kullback-Leibler information function to the sequential design of experiments when the cost of experimenting is small. His results are valid for the case of two terminal decisions and a finite number of experiments and states of nature. These results have been generalized by Albert (1961) to an infinite number of states of nature and by Bessler (1960) to an infinite number of experiments and k terminal actions. Kiefer and Sacks (1963) have also provided some extensions.

The statement of Chernoff's problem and the problem considered here are not identical and we proceed by analogizing his results to the problem at hand. In the context of the current problem, he would proceed by first assuming that at each stage of sampling the model with the largest posterior probability is the correct one. Then if A denotes the space of allowable experiments, define the Kullback-Leibler (K-L) information about model j in experiment $a \in A$ when model i is true as

$$I(a, i, j) = \int \ln \left[\frac{f_i(y|a)}{f_j(y|a)} \right] f_i(y|a) dy$$

where $f_i(y|a)$ denotes the probability density of y under model i when experiment $a \in A$ is performed. Let \hat{i} denote the model with the highest current probability of being the correct one. Then in analogy to Chernoff, we define the optimal experiment as $a(\hat{i})$ where $a(\hat{i})$ is defined by that experiment satisfying

$$I[a(\hat{i}), \hat{i}, j] = \sup_{a \in A} \inf_{k \neq \hat{i}} I(a, \hat{i}, k)$$

That is, Chernoff represents the problem as a game between nature and the statistician where the statistician maximizes over A and nature minimizes over the alternative models assuming \hat{i} is the correct model. Chernoff also specifically derives a stopping rule which we do not discuss here.

Hunter and Reiner (1965) considered a sequential design procedure for discriminating between two model equations. Their procedure chooses the experimental conditions which, based upon maximum likelihood estimates of the parameters from the data already collected, separate the expected values of the observed variable under the two models by as much as possible.

Box and Hill (1967) discussed the use of the Kullback-Leibler information function, deriving it from considerations involving the entropy function. They consider the use of the K-L information function to sequentially discriminate among several mechanistic (nonlinear) model equations. Besides the fact that they consider

nonlinear models, their approach is different in the sense that although they do assume prior probabilities on the proposed models, and compute posterior probabilities from the observations, they assume the parameters of the model equations are known constants.

Meeter, Pirie, and Blot (1970) have done a number of computer simulations comparing the methods of Chernoff and of Box and Hill. They found that the Box-Hill procedure performed quite well on the examples in comparison to Chernoff's procedure. It is interesting to note that Chernoff seems to be the only one of these authors who defined an explicit rule for terminating sampling. Although Chernoff's procedure is known to be asymptotically optimal, it is also known to require very large sample sizes.

CHAPTER 2

STRUCTURE OF THE LINEAR MODELS

In the theory of the general linear statistical model, we are concerned with problems involving model equations relating k controlled variables (z_i ; $i = 1, \dots, k$) to an observed variable (y). The form of the model equation is required to be

$$y = \sum_{i=1}^I \beta_i h_i(z_1, \dots, z_k) + \epsilon$$

The known functions h_i are arbitrary except that they may not contain any unknown parameters. The equation is linear in the unknown parameters β_i and ϵ is assumed to be a random variable with expectation zero and known finite variance. We may write

$x_i = h_i(z_1, \dots, z_k)$ and henceforth express the models in terms of the x_i variables. If n observations are made upon y we let x_{ji} denote the value of x_i at which the j^{th} observation is made. Thus for the n observations the model may conveniently be written as

$$\vec{y} = M\vec{\beta} + \vec{\epsilon} \quad (2-1)$$

where

$$\vec{y}' = (y_1, y_2, \dots, y_n)$$

$$M = \begin{bmatrix} x_{11} & x_{12} & \dots & x_{1I} \\ x_{21} & x_{22} & & x_{2I} \\ \cdot & \cdot & & \cdot \\ \cdot & \cdot & & \cdot \\ \cdot & \cdot & & \cdot \\ x_{n1} & x_{n2} & \dots & x_{nI} \end{bmatrix}$$

$$\vec{\beta}' = (\beta_1, \beta_2, \dots, \beta_I)$$

$$\vec{\epsilon}' = (\epsilon_1, \epsilon_2, \dots, \epsilon_n)$$

and the ϵ_i are uncorrelated. The matrix M is called the design matrix for the experiment consisting of the n observations. The problem of experimental design is that of choosing the x_{ji} values in some "optimal" manner.

In certain situations in practice the experimenter can postulate several possible models involving different functions of the z_i variables which correspond to several possible mechanistic or empirically based theories. These may lead to the various models containing different sets of x_i . There may be some overlapping of the x_i among the models or there may be none.

There are then two problems requiring solution. The first is that of choosing experiment designs which will enable the experimenter to decide which of the potential models is the correct one. Then, having chosen the model, the parameters must be estimated. The second problem has many solutions using a variety of standard

techniques. This dissertation concerns itself primarily with a method of designing experiments to provide information for choosing the appropriate model equation.

We assume there are L different competing model equations. These models may be combined into one large possible model equation and then the L hypothetical models are equivalent to there being L hypotheses restricting certain sets of parameters of the large model to be a priori zero. For example, we might have two controlled variables x_1 and x_2 . And suppose the model equations postulated are:

$$H_1: y = \beta_1^{(1)} x_1 + \epsilon$$

$$H_2: y = \beta_2^{(2)} x_2 + \epsilon$$

$$H_3: y = \beta_1^{(3)} x_1 + \beta_2^{(3)} x_2 + \epsilon$$

where $\beta_i^{(j)}$ denotes the coefficient of controlled variable i in model equation j . The distinction must be made because although $\beta_i^{(j)}$ and $\beta_i^{(k)}$ are coefficients of variable i , their distributions need not be the same. This notation is clumsy, however, and if we implicitly accept the fact that the distributions of the $\beta_i^{(j)}$ depend upon the model, we may more simply rewrite the models as

$$H_1: y = \beta_1 x_1 + \epsilon$$

$$H_2: y = \beta_2 x_2 + \epsilon$$

$$H_3: y = \beta_1 x_1 + \beta_2 x_2 + \epsilon$$

We say that models 1 and 2 are nested within model 3. This is

equivalent to writing one model as $y = \beta_1 x_1 + \beta_2 x_2 + \epsilon = \vec{X}'\vec{B} + \epsilon$ and hypothesizing

$$H_1: \beta_2 = 0$$

$$H_2: \beta_1 = 0$$

$$H_3: \beta_1 \neq 0, \beta_2 \neq 0$$

In this sense it is seen that the terms model and hypothesis are interchangeable and will be used interchangeably in the remainder of this dissertation. The notation we adopt is that H_ℓ claims

$$\vec{y} = M_\ell \vec{\alpha}_\ell + \vec{\epsilon}$$

where $\vec{\alpha}_\ell$ is the appropriate $k_\ell \times 1$ vector of β 's from \vec{B} which appear in model ℓ and M_ℓ is the appropriate matrix of x 's.

We now precisely state the three basic distributional assumptions about the parameters and random variables of the models:

(1) The vector $\vec{\epsilon}$ follows a multivariate normal distribution with mean $\vec{0}$ and precision matrix T . T is assumed known. (The precision matrix is the inverse of the covariance matrix of the distribution.) Since T must be positive definite symmetric, we need only consider the special case where $T = \tau I$ since linear transformation of the y reduces all other cases to this one. Note that we assume τ is known. Thus $\vec{\epsilon} \sim N(\vec{0}, \tau I)$.

(2) For each $\ell = 1, \dots, L$ the prior distribution of $\vec{\alpha}_\ell$ is

$$\vec{\alpha}_\ell \sim N(\vec{\mu}_{\ell,0}, \Psi_{\ell,0})$$

where $\vec{\mu}_{\ell,0}$ and $\Psi_{\ell,0}$ are known.

(3) The prior probability that the ℓ^{th} model is the correct

model equation is assumed specified and denoted by $\theta_{\ell,0}$. We require $\sum_{\ell=1}^L \theta_{\ell,0} = 1.0$. In order to satisfy this requirement in a completely precise manner we must make the models mutually exclusive. As described so far, this need not be true. However, this is a simple problem to get around for the following reason. Each of the H_{ℓ} specifies that $\vec{\alpha}_{\ell}$ is an element of a k_{ℓ} dimensional subset of K -space which we denote as E_{ℓ} . For any pair H_j and H_k we have either (1) $E_j \supset E_k$, (2) $E_k \supset E_j$, or (3) neither space contains the other and $E_j \cap E_k$ has measure zero with respect to H_j and H_k . For case 1 define θ_j as $\theta_j = \Pr\{\vec{\alpha}_j \in E_j - E_k\}$. But E_k has zero measure with respect to H_j and hence the distribution function of $\vec{\alpha}_j$ restricted to $E_j - E_k$ is identical to the distribution function of $\vec{\alpha}_j$ over all of E_j . Thus for any practical purpose, the fact that $E_j \supset E_k$ does not affect any probability computations. Similar arguments apply to cases two and three.

We now describe the space A of allowable experiments in more detail. If the number of elements of \vec{X} is K , then a choice of experiment $a \in A$ is composed of the number J of observations to take and J vectors from some subset of Euclidean K -space. The J vectors specify the values of the controlled variables x_{ji} . At the j^{th} experiment or j^{th} stage of experimenting the particular choice from A is denoted a_j .

CHAPTER 3

PREREQUISITE DISTRIBUTION THEORY

In the remainder of this dissertation, much use will be made of the distribution of the observed variable, the posterior probabilities of the models, and the posterior distributions of the parameters of the model equations. The first part of this chapter develops these distributions. The second part derives the fact that the sampling procedure is Markovian in nature and provides a notation for describing the state of the process. The third section of this chapter discusses some results on the limiting behavior of the posterior distributions when the sequence of experiments is chosen in advance. The strong restrictions that must be made to accomplish these large sample results and the fact that they do not describe the adaptive process might lead one to believe that they are not worthwhile pursuing.

We find in chapter 6, however, that there is a close analogy between these results and the behavior of the adaptive procedure, and that these results help explain and interpret the Monte Carlo simulation results.

3.1 Posterior and Marginal Distributions

Let $f_{\ell}(\vec{y}_{j+1} | a_{j+1}, \vec{\alpha}_{\ell})$ denote the density function of the vector \vec{y}_{j+1} under H_{ℓ} when the parameter values are given by $\vec{\alpha}_{\ell}$ at stage $j + 1$ of sampling. Let the probability density function

of $\vec{\alpha}_\ell$ after j stages of sampling be denoted $\xi_{\ell,j}(\vec{\alpha})$. This is a preposterior density since it serves as the posterior density of $\vec{\alpha}_\ell$ after j stages of sampling and the prior density of $\vec{\alpha}_\ell$ before the $j+1^{\text{st}}$ stage of sampling occurs.

Lemma 3.1: After j stages of sampling, $\vec{\alpha}_\ell$ follows a multivariate normal distribution with mean vector $\vec{\mu}_{\ell,j}$ and precision matrix

$\Psi_{\ell,j}$. That is, after j stages of sampling,

$$\vec{\alpha}_\ell \sim N(\vec{\mu}_{\ell,j}, \Psi_{\ell,j})$$

where

$$\begin{aligned} \Psi_{\ell,j} &= \Psi_{\ell,j-1} + M'_{\ell,j} TM_{\ell,j} \\ &= \Psi_{\ell,0} + \sum_{i=1}^j M'_{\ell,i} TM_{\ell,i} \end{aligned} \quad (3-1)$$

and

$$\begin{aligned} \vec{\mu}_{\ell,j} &= \Psi_{\ell,j}^{-1} (M'_{\ell,j} T \vec{y}_j + \Psi_{\ell,j-1} \vec{\mu}_{\ell,j-1}) \\ &= \Psi_{\ell,j}^{-1} \left[\sum_{i=1}^j M'_{\ell,i} T \vec{y}_i + \Psi_{\ell,0} \vec{\mu}_{\ell,0} \right] \end{aligned} \quad (3-2)$$

and where $M_{\ell,i}$ denotes the design matrix specified by a_i under H_ℓ .

Proof: By Bayes theorem, if \vec{y}_j is the observed vector at stage j

$$\begin{aligned} \xi_{\ell,j}(\vec{\alpha}) &= \frac{f_\ell(\vec{y}_j | a_j, \vec{\alpha}) \xi_{\ell,j-1}(\vec{\alpha})}{\int f_\ell(\vec{y}_j | a_j, \vec{\alpha}^*) \xi_{\ell,j-1}(\vec{\alpha}^*) d\vec{\alpha}^*} \\ &\propto f_\ell(\vec{y}_j | a_j, \vec{\alpha}) \xi_{\ell,j-1}(\vec{\alpha}) \end{aligned} \quad (3-3)$$

The symbol \propto means "proportional to" and is used in the context of DeGroot's (1970, p. 160) usage. Thus

$$\xi_{\ell,j}(\vec{\alpha}) \propto e^{-Q/2}$$

where (dropping subscripts)

$$Q = (\vec{M}\vec{\alpha} - \vec{y})' T (\vec{M}\vec{\alpha} - \vec{y}) + (\vec{\alpha} - \vec{\mu})' \Psi (\vec{\alpha} - \vec{\mu})$$

Since T and Ψ are positive definite symmetric we can write

$$\begin{aligned} Q &= \vec{\alpha}' (\Psi + M' T M) \vec{\alpha} - 2(\vec{\mu}' \Psi + \vec{y}' T M) \vec{\alpha} + (\text{terms not involving } \vec{\alpha}) \\ &= [\vec{\alpha} - (\Psi + M' T M)^{-1} (M' T \vec{y} + \Psi \vec{\mu})]' \cdot (\Psi + M' T M) \\ &\quad \cdot [\vec{\alpha} - (\Psi + M' T M)^{-1} (M' T \vec{y} + \Psi \vec{\mu})] + (\text{terms not involving } \vec{\alpha}) \end{aligned}$$

The terms not involving $\vec{\alpha}$ may be factored out through use of the proportionality device, leaving the kernel of a multivariate normal distribution with parameters as specified by the lemma. Thus $\vec{\alpha}_{\ell}$ is distributed as claimed.

Q.E.D.

Owen (1970) has derived a result similar to Lemma 3.1 in the case of a two factor experiment where the factors are treatments and blocks.

We now turn to determining the distribution of \vec{y}_{j+1} . This is done in two stages. First we do not know which of the models is in fact the correct one. Then for any given model, we do not know the value of $\vec{\alpha}_{\ell}$. Let $f_{\ell}(\vec{y}_{j+1} | a_{j+1}, \vec{\alpha})$ denote the distribution of \vec{y}_{j+1} under H_{ℓ} when experiment $a_{j+1} \in A$ is performed and $\vec{\alpha}_{\ell}$ is specified. Since we do not know $\vec{\alpha}_{\ell}$ we must average this distribution over all $\vec{\alpha}_{\ell}$. Let $f_{\ell}(\vec{y}_{j+1} | a_{j+1})$ denote the mixture of the densities $f_{\ell}(\vec{y}_{j+1} | a_{j+1}, \vec{\alpha})$ with respect to the marginal posterior of $\vec{\alpha}_{\ell}$.

Lemma 3.2 The conditional distribution of \vec{y}_j given H_{ℓ} and a_j is a multivariate normal distribution with mean vector $\vec{s}_{\ell,j}$ and precision matrix $R_{\ell,j}$ where

$$R_{\ell,j} = T \left[I - M_{\ell,j} (M'_{\ell,j} T M_{\ell,j} + \Psi_{\ell,j-1})^{-1} M'_{\ell,j} T \right] \quad (3-4)$$

$$\vec{s}_{\ell,j} = R_{\ell,j}^{-1} T M_{\ell,j} (M'_{\ell,j} T M_{\ell,j} + \Psi_{\ell,j-1})^{-1} \Psi_{\ell,j-1} \vec{\mu}_{\ell,j-1} \quad (3-5)$$

Proof: The required mixture distribution is given by

$$\begin{aligned} f_{\ell}(\vec{y}_j | a_j) &= \int f_{\ell}(\vec{y}_j | a_j, \vec{\alpha}) \pi_{\ell,j-1}(\vec{\alpha}) d\vec{\alpha} \\ &= \int e^{-Q/2} d\vec{\alpha} \end{aligned}$$

where

$$\begin{aligned} Q &= (\vec{y}_j - M_{\ell,j} \vec{\alpha})' T (\vec{y}_j - M_{\ell,j} \vec{\alpha}) + (\vec{\alpha} - \vec{\mu}_{\ell,j-1})' \Psi_{\ell,j-1} (\vec{\alpha} - \vec{\mu}_{\ell,j-1}) \\ &= \vec{\alpha}' (M' T M + \Psi) \vec{\alpha} - 2 \vec{\alpha}' (M' T M + \Psi) (M' T M + \Psi)^{-1} (M' T \vec{y} + \Psi \vec{\mu}) \\ &\quad + \vec{y}' T \vec{y} + \vec{\mu}' \Psi \vec{\mu} \\ &= \vec{\alpha}' (M' T M + \Psi) \vec{\alpha} - 2 \vec{\alpha}' (M' T M + \Psi) (M' T M + \Psi)^{-1} (M' T \vec{y} + \Psi \vec{\mu}) \\ &\quad + (M' T \vec{y} + \Psi \vec{\mu})' (M' T M + \Psi)^{-1} (M' T \vec{y} + \Psi \vec{\mu}) \\ &\quad - (M' T \vec{y} + \Psi \vec{\mu})' (M' T M + \Psi)^{-1} (M' T \vec{y} + \Psi \vec{\mu}) + \vec{y}' T \vec{y} + \vec{\mu}' \Psi \vec{\mu} \end{aligned}$$

The first three terms yield the quadratic form

$$\begin{aligned} Q_1 &= [\vec{\alpha} - (M' T M + \Psi)^{-1} (M' T \vec{y} + \Psi \vec{\mu})]' \cdot (M' T M + \Psi) \\ &\quad \cdot [\vec{\alpha} - (M' T M + \Psi)^{-1} (M' T \vec{y} + \Psi \vec{\mu})] \end{aligned}$$

The remainder of Q does not involve $\vec{\alpha}$ and $e^{-Q_1/2}$ is the kernel of a multivariate normal distribution so that when $e^{-Q_1/2}$

is integrated over $\vec{\alpha}$ we remain with

$$f_{\ell,j-1}(\vec{y}_j | a_j) \propto e^{-Q_2/2}$$

where

$$\begin{aligned} Q_2 &= -(M'T\vec{y} + \vec{\psi}\vec{\mu})'(M'TM + \Psi)^{-1}(M'T\vec{y} + \vec{\psi}\vec{\mu}) + \vec{y}'T\vec{y} + \vec{\mu}'\Psi\vec{\mu} \\ &= y'[T - TM(M'TM + \Psi)^{-1}M'T]\vec{y} - 2\vec{y}'TM(M'TM + \Psi)^{-1}\vec{\psi}\vec{\mu} \\ &\quad + (\text{terms not involving } \vec{y}) \\ &= (\vec{y} - \vec{s})'R(\vec{y} - \vec{s}) + (\text{terms not involving } \vec{y}) \end{aligned}$$

The terms not involving \vec{y} may be factored out via the proportionality device leaving

$$f_{\ell}(\vec{y}_j | a_j) \propto e^{-(\vec{y} - \vec{s}_{\ell,j})'R_{\ell,j}(\vec{y} - \vec{s}_{\ell,j})/2}$$

This is the kernel of a multivariate normal distribution with mean vector $\vec{s}_{\ell,j}$ and precision matrix $R_{\ell,j}$ as claimed. Thus the density of y_j given H_{ℓ} and a_j is given by

$$f_{\ell}(\vec{y}_j | a_j) = (2\pi)^{-J/2} |R_{\ell,j}|^{1/2} \exp\left\{-\frac{1}{2} (\vec{y}_j - \vec{s}_{\ell,j})'R_{\ell,j}(\vec{y}_j - \vec{s}_{\ell,j})\right\} \quad (3-6)$$

Q.E.D.

Since the true model is unknown we now compute the mixture of the distributions of Lemma 3.2 with respect to the probabilities $\theta_{\ell,j}$ as

$$f(y_j | a_j) = \sum_{\ell=1}^L \theta_{\ell,j-1} f_{\ell}(\vec{y}_j | a_j) \quad (3-7)$$

To compute the posterior probability of each model being cor-

rect after the observation y_{j+1} is obtained, we apply Bayes theorem directly to get

$$\theta_{\ell,j+1} = \frac{f_{\ell}(\vec{y}_{j+1} | a_{j+1}) \theta_{\ell,j}}{\sum_{k=1}^L f_k(\vec{y}_{j+1} | a_{j+1}) \theta_{k,j}} \quad (3-8)$$

3.2 Markovian Nature of Sampling Process

Consider a sequence of random variables W_1, W_2, \dots which take on values in a sample space or state space Ω . We let \mathcal{F} denote the σ -field of subsets of Ω for which probabilities are defined. The sequence of random variables W_i form a Markov Process if for every $F \in \mathcal{F}$ and for all w_1, \dots, w_n in Ω , and all for n , $n = 1, 2, 3, \dots$ we have

$$\begin{aligned} \Pr\{W_{n+1} \in F | W_1 = w_1, \dots, W_n = w_n\} \\ &= \Pr\{W_{n+1} \in F | W_n = w_n\} \\ &= \int_F g_{n+1}(w | w_n) dw \end{aligned} \quad (3-9)$$

where $g_{n+1}(w | w_n)$ denotes the generalized conditional probability density function of W_{n+1} . If the conditional probabilities in equation (3-9), equivalently the g_n , do not depend upon n the transition process is called stationary. The state space in this paper can be described by a vector containing: (1) the probabilities θ_{ℓ} , (2) the elements of the vectors describing the current posterior means under the various H_{ℓ} , and (3) the lower triangular part of the current posterior precision matrices under the various H_{ℓ} . Thus

$$\Omega = \left\{ \theta_1, \dots, \theta_L, \mu_1^{(1)}, \mu_1^{(2)}, \dots, \mu_2^{(1)}, \mu_2^{(2)}, \dots, \mu_L^{(1)}, \dots, \right. \\ \left. \mu_L^{(K_L)}, \psi_1^{(1,1)}, \dots, \psi_L^{(K_L, K_L)} : \sum_{i=1}^L \theta_i = 1.0, -\infty < \mu_k^{(j)} < +\infty, \right. \\ \left. \psi_i \text{ positive definite for } i = 1, \dots, L \right\} \quad (3-10)$$

For any given state $w \in \Omega$ the transition to the next state depends only upon the state w and the experiment $a \in A$ that is chosen.

This is true because a determines the posterior precision matrices regardless of the value of y , and the posterior means $\vec{\mu}_g(y)$ and probabilities $\vec{\theta}(y)$ are determined by equations (3-1) and (3-8) which again depend only upon w , y , and a . Thus the transition process on the states is Markovian. The process is stationary, also, since for given initial g_1 the successive g_n do not depend upon n . The transition function may be described as follows. Define a mapping $T: \Omega \times Y \rightarrow \Omega$ and let $Q(\Omega \times Y)$ denote the Borel sets on $\Omega \times Y$ and $Q(\Omega)$ denote the Borel sets on Ω . Let $T^{-1}(F)$ denote the inverse image of F where $F \in Q(\Omega)$ and $T^{-1}(F) \in Q(\Omega \times Y)$. Then if w' denotes the state of the system after sampling,

$$\Pr(w' \in F | w, a) = \int \sum_{i=1}^L \theta_i f_i(y | a, w) dy \\ (w, y) \in T^{-1}(F)$$

3.3 Large Sample and Limiting Results

Even though this paper is concerned primarily with small sample procedures, it is interesting and informative to know the large sample or limiting behavior of the parameters and the sampling

process. Unfortunately, for the adaptive procedure this is an extremely difficult subject to study. Thus we do not study the adaptive procedure here but instead consider the experiment selection procedure under the restrictions listed below in the hope that these results will illuminate the adaptive procedure in some sense.

(1) Assume A is finite with $N(A)$ elements, and represented as

$$A = \{a^{(1)}, a^{(2)}, \dots, a^{[N(A)]}\}$$

(2) An infinite sequence $\{a_j\}$ is specified such that as the number of experiments approaches infinity, the proportion of times that $a^{(i)}$ is performed approaches p_i with $0 < p_i < 1$ and $\sum p_i = 1.0$. The experiments a_j are chosen independently of each other.

(3) Assume H_{i*} is the true model and that $\vec{\mu}^*$ is the true value of the parameters in the model.

(4) Assume that only one observation is taken in each experiment $a^{(i)}$.

(5) Assume that the structure of A is such that all matrices under consideration are nonsingular.

It should be noted that the most restrictive of the above assumptions is the second. For in a true sequential decision procedure, the actual experiment chosen is a random variable depending upon the previous observations obtained. Since we are in fact studying a problem other than the one of most importance the remainder of the chapter will not be developed in rigorous detail and the results

obtained cannot be rigorously applied to the sequential procedure. It will be seen in Chapter 6, however, that fairly extensive Monte Carlo simulations seem to bear up the general conclusions reached here.

Let $k(j)$ denote the superscript of the experiment performed at stage j . Thus if $a_{10} = a^{(5)}$, then $k(10) = 5$. Also let $n(i, j)$ denote the number of times $a^{(i)}$ is performed in the sequence of experiments up to and including the j^{th} stage. Let $M_{\ell, i}$ denote the design matrix under H_{ℓ} when $a^{(i)}$ is chosen.

Lemma 3.3 Under the above assumptions the posterior precision matrices and mean vectors converge with probability one as $j \rightarrow \infty$ to:

$$\frac{1}{j} \tau \Psi_{\ell, j} \rightarrow \Psi_{\ell} = \sum_{i=1}^{N(A)} p_i M'_{\ell, i} M_{\ell, i}$$

$$\vec{\mu}_{\ell, j} \rightarrow (\Psi_{\ell})^{-1} \left(\sum_{i=1}^{N(A)} p_i M'_{\ell, i} M_{\ell, i}^* \right) \vec{\mu}^*$$

Proof: To prove the first limit, recall from equation (3-1) that

$$\Psi_{\ell, j} = \Psi_{\ell, 0} + \sum_{i=1}^j \tau M'_{\ell, k(i)} M_{\ell, k(i)}$$

Thus

$$\begin{aligned}\frac{1}{j\tau} \psi_{\ell,j} &= \frac{1}{j\tau} \psi_{\ell,0} + \sum_{i=1}^j \frac{1}{j} M'_{\ell,k(i)} M_{\ell,k(i)} \\ &= \frac{1}{j\tau} \psi_{\ell,0} + \sum_{i=1}^{N(A)} \frac{n(i,j)}{j} M'_{\ell,i} M_{\ell,i}\end{aligned}$$

As $j \rightarrow \infty$, the first term goes to zero and the factors $\frac{n(i,j)}{j} \rightarrow p_i$ by assumption.

To prove the second part note that $y_j = M_{i^*,k(j)} \vec{\mu}^* + \epsilon_j$ where $\epsilon_j \sim N(0, \tau)$. Using the second form of equation (3-2) we get

$$\vec{\mu}_{\ell,j} = (\psi_{\ell,j})^{-1} \left\{ \tau \sum_{i=1}^j M'_{\ell,k(i)} y_i + \psi_{\ell,0} \vec{\mu}_{\ell,0} \right\} \quad (3-11)$$

Then substituting the expression for y_j into equation (3-11) gives

$$\begin{aligned}\vec{\mu}_{\ell,j} &= (\psi_{\ell,j})^{-1} \left\{ \sum_{i=1}^{N(A)} n(i,j) \tau M'_{\ell,i} M_{i^*,i} \vec{\mu}^* \right. \\ &\quad \left. + \tau \sum_{i=1}^{N(A)} M'_{\ell,i} \left(\sum_{m=1}^j \epsilon_m \delta_{k(m),i} \right) + \psi_{\ell,0} \vec{\mu}_{\ell,0} \right\}\end{aligned}$$

where $\delta_{i,j}$ denotes the Kronecker delta function. Thus

$$\begin{aligned}\vec{\mu}_{\ell,j} &= \left(\frac{1}{j\tau} \psi_{\ell,j} \right)^{-1} \left\{ \sum_{i=1}^{N(A)} \frac{n(i,j)}{j} M'_{\ell,i} M_{i^*,i} \vec{\mu}^* \right. \\ &\quad \left. + \sum_{i=1}^{N(A)} M'_{\ell,i} \left(\frac{1}{j} \sum_{m=1}^j \epsilon_m \delta_{k(m),i} \right) + \frac{1}{j\tau} \psi_{\ell,0} \vec{\mu}_{\ell,0} \right\}\end{aligned}$$

From assumption 2 we know that $n(i,j) \rightarrow \infty$ as $j \rightarrow \infty$ and since the

ε_m form a sequence of independent and identically distributed random variables, the strong law of large numbers may be applied to show for $i = 1, \dots, L$

$$\Pr \left\{ \lim_{n(i,j) \rightarrow \infty} \frac{1}{n(i,j)} \sum_{m=1}^j \varepsilon_m \delta_{k(m),i} = 0 \right\} = 1.0$$

Since $\frac{1}{j\tau} \rightarrow 0$ as $j \rightarrow \infty$ we then have

$$\vec{\mu}_{\ell,j} \approx \Psi_{\ell}^{-1} \left(\sum_{i=1}^{N(A)} \frac{n(i,j)}{j} M'_{\ell,i} M_{i^*,i} \right)^{\rightarrow*}_{\mu}$$

This sequence will not have a limit unless $\lim_{j \rightarrow \infty} \frac{n(i,j)}{j} = p_i$ exists. If such a limit exists, the lemma follows immediately.

Q.E.D.

Lemma 3.4 Under the assumptions stated, $R_{\ell,j} \rightarrow \tau$ irrespective of $\{a_j\}$ and

$$s_{\ell,j} \approx M_{\ell,k(j)} (\Psi_{\ell})^{-1} \left(\sum_{i=1}^{N(A)} p_i M'_{\ell,i} M_{i^*,i} \right)^{\rightarrow*}_{\mu}$$

for large enough j .

Proof: From equation (3-4) and the assumptions

$$R_{\ell,j} = \tau \left\{ 1.0 - \tau M_{\ell,k(j)} [\Psi_{\ell,j} + \tau M'_{\ell,k(j)} M_{\ell,k(j)}]^{-1} M'_{\ell,k(j)} \right\}$$

As $j \rightarrow \infty$, $[\Psi_{\ell,j} + \tau M'_{\ell,k(j)} M_{\ell,k(j)}] \rightarrow (\infty)$ and hence its inverse $\rightarrow (0)$. But then $R_{\ell,j} \rightarrow \tau$ as claimed.

From equation (3-5)

$$s_{\ell,j} = (R_{\ell,j})^{-1} \left\{ \tau M_{\ell,k(j)} [M'_{\ell,k(j)} M_{\ell,k(j)} \tau + \Psi_{\ell,j}]^{-1} \Psi_{\ell,j} \vec{\mu}_{\ell,j} \right\}$$

For large enough j , $[M'_{\ell,k(j)} M_{\ell,k(j)}^\tau + \Psi_{\ell,j}]^{-1} \Psi_{\ell,j}$ is asymptotically like the identity matrix, I , so that

$$\begin{aligned} s_{\ell,j} &\approx \frac{1}{\tau} \left\{ \tau M_{\ell,k(j)} \vec{\mu}_{\ell,j} \right\} \\ &\approx M_{\ell,k(j)} (\Psi_{\ell})^{-1} \left[\sum_{i=1}^{N(A)} p_i M'_{\ell,i} M_{i^*,i} \right] \vec{\mu}^* \end{aligned}$$

We note that if $\ell = i^*$, then from the definition of Ψ_{ℓ} , we have $s_{i^*,j} \approx M_{i^*,k(j)} \vec{\mu}^*$ as expected.

Q.E.D.

Lemma 3.5 If H_{i^*} is the true hypothesis and the model of H_{i^*} is nested within the model of H_{ℓ} , then under the above assumptions and assuming the parameter vectors are rearranged appropriately

$$\Pr \left\{ \vec{\mu}_{\ell,j} \rightarrow \begin{pmatrix} \vec{\mu}^* \\ \vec{0} \end{pmatrix} \right\} = 1.0$$

Proof: From Lemma 3.3 we have

$$\vec{\mu}_{\ell,j} \rightarrow \Psi_{\ell}^{-1} \left(\sum_{i=1}^{N(A)} p_i M'_{\ell,i} M_{i^*,i} \right) \vec{\mu}^*$$

If the parameters are rearranged appropriately then $M_{\ell,i}$ may be written

$$M_{\ell,i} = (M_{i^*,i}, \tilde{M}_{\ell,i})$$

where $\tilde{M}_{\ell,i}$ denotes the design matrix corresponding to the independent variables in H_{ℓ} but not in H_{i^*} . Thus

$$M'_{\ell,i} M_{\ell,i} = \begin{bmatrix} M'_{i^*,i} M_{i^*,i} & M'_{i^*,i} \tilde{M}_{\ell,i} \\ \tilde{M}'_{\ell,i} M_{i^*,i} & \tilde{M}'_{\ell,i} \tilde{M}_{\ell,i} \end{bmatrix}$$

and

$$\Psi_{\ell} = \begin{bmatrix} \Psi_{i^*} & \Psi_{i^*,\ell} \\ \Psi_{i^*,\ell} & \tilde{\Psi}_{\ell} \end{bmatrix}$$

where

$$\Psi_{i^*} = \sum_{i=1}^{N(A)} p_{iM'_{i^*},i} M_{i^*,i}$$

$$\Psi_{i^*,\ell} = \sum_{i=1}^{N(A)} p_{iM'_{i^*},i} \tilde{M}_{\ell,i}$$

and

$$\tilde{\Psi}_{\ell} = \sum_{i=1}^{N(A)} p_{i\tilde{M}_{\ell},i} \tilde{M}_{\ell,i}$$

Thus from a well known identity (e.g., Graybill (1969), p. 165)

$$\Psi_{\ell}^{-1} = \begin{bmatrix} (\Psi_{i^*} - \Psi_{i^*,\ell} \tilde{\Psi}_{\ell}^{-1} \Psi_{i^*,\ell})^{-1} & -\Psi_{i^*,\ell}^{-1} \tilde{\Psi}_{\ell}^{-1} (\tilde{\Psi}_{\ell} - \Psi_{i^*,\ell} \Psi_{i^*,\ell}^{-1} \Psi_{i^*,\ell}) \\ -(\tilde{\Psi}_{\ell} - \Psi_{i^*,\ell} \Psi_{i^*,\ell}^{-1} \Psi_{i^*,\ell})^{-1} \Psi_{i^*,\ell}^{-1} & (\tilde{\Psi}_{\ell} - \Psi_{i^*,\ell} \Psi_{i^*,\ell}^{-1} \Psi_{i^*,\ell}) \end{bmatrix}$$

Also

$$\sum_{i=1}^{N(A)} p_{iM'_{\ell},i} M_{i^*,i} = \sum_{i=1}^{N(A)} \begin{pmatrix} p_{iM'_{i^*},i} M_{i^*,i} \\ p_{i\tilde{M}_{\ell},i} M_{i^*,i} \end{pmatrix} = \begin{pmatrix} \Psi_{i^*} \\ \Psi_{i^*,\ell} \end{pmatrix}$$

Thus

$$\Psi_{\ell}^{-1} \begin{pmatrix} \Psi_{i^*} \\ \Psi_{i^*, \ell} \end{pmatrix} = \begin{bmatrix} \left(\Psi_{i^*} - \Psi_{i^*, \ell} \tilde{\Psi}_{\ell}^{-1} \Psi_{i^*, \ell}' \right)^{-1} \Psi_{i^*} \\ - \Psi_{i^*, \ell}' \Psi_{i^*, \ell} \left(\tilde{\Psi}_{\ell} - \Psi_{i^*, \ell}' \Psi_{i^*, \ell}^{-1} \Psi_{i^*, \ell} \right)^{-1} \Psi_{i^*, \ell}' \\ - \left(\tilde{\Psi}_{\ell} - \Psi_{i^*, \ell}' \Psi_{i^*, \ell}^{-1} \Psi_{i^*, \ell} \right)^{-1} \Psi_{i^*, \ell} \\ + \left(\tilde{\Psi}_{\ell} - \Psi_{i^*, \ell}' \Psi_{i^*, \ell}^{-1} \Psi_{i^*, \ell} \right)^{-1} \Psi_{i^*, \ell} \end{bmatrix}$$

$$= \begin{bmatrix} I \\ 0 \end{bmatrix}$$

upon application of Lemma 3.6 which follows. Thus

$$\vec{\mu}_{\ell, j} \rightarrow \begin{pmatrix} I \\ 0 \end{pmatrix} \vec{\mu}^* = \begin{pmatrix} \vec{\mu}^* \\ 0 \end{pmatrix}$$

Q.E.D.

Lemma 3.6. (Problem 2.9 of Rao (1965)) If A and D are matrices possessing inverses, then

$$(A + BDB')^{-1} = A^{-1} - A^{-1}B(B'A^{-1}B + D^{-1})^{-1}B'A^{-1}$$

Proof: By direct multiplication we only need show

$$\begin{aligned} I &= (A + BDB')(A^{-1} - A^{-1}B(B'A^{-1}B + D^{-1})^{-1}B'A^{-1}) \\ &= I - B(B'A^{-1}B + D^{-1})^{-1}B'A^{-1} + BDB'A^{-1} \\ &\quad - BDB'A^{-1}B(B'A^{-1}B + D^{-1})^{-1}B'A^{-1} \\ &= I - B[-(B'A^{-1}B + D^{-1})^{-1} + D - DB'A^{-1}B(B'A^{-1}B + D^{-1})^{-1}]B'A^{-1} \\ &= I - B[D - [I + DB'A^{-1}B][B'A^{-1}B + D^{-1}]^{-1}]B'A^{-1} \\ &= I - B[D - D[D^{-1} + B'A^{-1}B][D^{-1} + B'A^{-1}B]^{-1}]B'A^{-1} \\ &= I - B[D - D]B'A^{-1} \\ &= I \end{aligned}$$

Q.E.D.

To apply this result to Lemma 3.5 simply set

$$A = \Psi_{i*}$$

$$B = \Psi_{i*,\ell}$$

$$D = -\tilde{\Psi}_{\ell}^{-1}$$

We now turn to consideration of the limiting behavior of $\theta_{\ell,j}$.

Computer simulations for both nested and non-nested cases indicate that for any k where H_{i*} is not nested in H_k , $\theta_{k,j} \rightarrow 0.0$ fairly rapidly and steadily. If H_{i*} is nested in H_k then it seems that $\theta_{k,j} \rightarrow 0.0$. The rate is initially rapid but then becomes very slow and it behaves in a very erratic manner. These points are discussed in some detail in Chapter 6.

It should be reiterated and these discussions have assumed the sequence $\{a_j\}$ to be specified and fixed for the sequence of experiments. In a sequential decision problem the sequence $\{a_j\}$ is not fixed, but $k(j)$ is in fact a random variable whose distribution depends upon $k(i)$ for $i < j$ and the y_i for $i < j$.

CHAPTER 4

ENTROPY FUNCTIONS AND THE KULLBACK-LEIBLER INFORMATION FUNCTION

When comparing a number of experiments to determine which is the optimal one to perform, one must define optimal. In this dissertation, that experiment which yields the largest expected K-L information is defined as the optimal experiment. In particular, let $I(w,a)$ denote the expected K-L information as a function of the experiment a and the current state w of the process. This function will be specified explicitly later. In this chapter, we first describe how the K-L information arises from attempting to reduce the entropy of the probabilities of the models. We then develop an expression for $I(w,a)$ and finally discuss the operational meaning of the use of $I(w,a)$ from a heuristic point of view.

4.1 Development of the K-L Information Function

The problem under consideration here is that we must choose one of a set of postulated model equations. For each model we have the posterior probability $\theta_{\ell,j}$ that it is the correct one. We would like to choose experiments which cause the posterior probability of the correct model to increase most rapidly. An indirect method of accomplishing this is to choose experiments which most rapidly decrease the entropy of the set of probabilities $\theta_{\ell,j}$. The entropy is defined as

$$\mathcal{E}(w) = - \sum_{\ell=1}^L \theta_{\ell,j} \ln(\theta_{\ell,j})$$

It can be verified that the entropy attains a maximum when all the probabilities are equal and attains a minimum when any one of the probabilities is one and the rest are zero.

Box and Hill (1967) proposed the use of the expected decrease between the entropy at the current stage of sampling and the anticipated entropy at the next stage of sampling as the criterion for selection of experiments. They found, however, that the entropy function is quite intractable analytically and applied a well-known inequality to show the expected K-L information function provides an upper bound on the reduction of entropy. Let $\theta_i(\vec{y}|w,a)$ denote the posterior probability of model i if the value \vec{y} is observed when the state was w . Let $w(\vec{y})$ denote the state of the process after observing the value \vec{y} when it was in state w . Then the anticipated entropy is given by

$$E\{\mathcal{E}[w(\vec{y}),a]\} = - \int \left\{ \sum_{\ell=1}^L \theta_{\ell}(\vec{y}|w,a) \ln[\theta_{\ell}(\vec{y}|w,a)] \right\} f(\vec{y}|w,a) d\vec{y}$$

Thus if the current state of the sampling process is $w \in \Omega$, and the experiment $a \in A$ is performed, the expected decrease in entropy, $R(w,a)$, is then defined as

$$\begin{aligned}
R(w, a) &= \mathcal{E}(w) - E\{\mathcal{E}[w(\vec{y}), a]\} \\
&= - \sum_{i=1}^L \theta_i \ln(\theta_i) + \int \left\{ \sum_{i=1}^L \theta_i(\vec{y}|w, a) \ln[\theta_i(\vec{y}|w, a)] \right\} \cdot \\
&\quad \left\{ \sum_{k=1}^L \theta_k f_k(\vec{y}|w, a) \right\} d\vec{y} \\
&= - \sum_{i=1}^L \theta_i \ln(\theta_i) + \int \sum_{\ell=1}^L \theta_\ell f_\ell(\vec{y}|w, a) \ln \left[\frac{\theta_\ell f_\ell(\vec{y}|w, a)}{\sum_{k=1}^L \theta_k f_k(\vec{y}|w, a)} \right] d\vec{y} \\
&\leq \int \sum_{\ell=1}^L \theta_\ell \left\{ \sum_{i=1}^L \theta_i f_i(\vec{y}|w, a) \ln \left[\frac{f_\ell(\vec{y}|w, a)}{f_i(\vec{y}|w, a)} \right] \right\} d\vec{y} \quad (4-1)
\end{aligned}$$

by application of the following inequality (Kullback (1969), p. 15)

$$\sum_{i=1}^L \theta_i f_i(\vec{y}|w, a) \ln \left[\frac{f_\ell(\vec{y}|w, a)}{f_i(\vec{y}|w, a)} \right] \geq f_\ell(\vec{y}|w, a) \ln \left[\frac{f_\ell(\vec{y}|w, a)}{\sum_{k=1}^L \theta_k f_k(\vec{y}|w, a)} \right]$$

Let

$$I(w, a, i, j) = \int f_i(\vec{y}|w, a) \ln \left[\frac{f_i(\vec{y}|w, a)}{f_j(\vec{y}|w, a)} \right] d\vec{y} \quad (4-2)$$

We note $I(w, a, i, j)$ is defined as the expected amount of information in the observations from experiment a for discriminating against H_j in favor of H_i . Let $Q(w, a)$ denote the matrix whose i, j element is $I(w, a, i, j)$. Then the inequality (4-1) may be written as

$$R(w,a) \leq \vec{\theta}' Q(w,a) \vec{\theta} = I(w,a) \quad (4-3)$$

Meeter et al. (1970) proposed the following heuristic argument in favor of using $I(w,a)$. If one knew that H_i were indeed the correct hypothesis and wished to maximize the information about H_k for $k \neq i$, then it would be natural to maximize

$$\sum_{k \neq i} \theta_k I(w,a,i,k)$$

But since H_i is assumed correct only with probability θ_i , it is equally natural to multiply the foregoing expression by θ_i and sum over i . But in doing this, one does end up with $I(w,a)$.

4.2 Evaluation of K-L Information Function

From equation (3-6) we have (if \vec{y} is $J \times 1$) that the density of \vec{y} under H_ℓ is given by

$$f_\ell(\vec{y}|a) = (2\pi)^{-J/2} |R_\ell|^{1/2} e^{-1/2(\vec{y}-\vec{s}_\ell)' R_\ell (\vec{y}-\vec{s}_\ell)}$$

Hence

$$\frac{f_m(\vec{y}|a)}{f_n(\vec{y}|a)} = \frac{|R_m|^{1/2} |R_n|^{-1/2} e^{-1/2(\vec{y}-\vec{s}_m)' R_m (\vec{y}-\vec{s}_m)}}{e^{-1/2(\vec{y}-\vec{s}_n)' R_n (\vec{y}-\vec{s}_n)}}$$

Moreover

$$\begin{aligned} \ln \left[\frac{f_m(\vec{y}|a)}{f_n(\vec{y}|a)} \right] &= \frac{1}{2} (\ln |R_m| - \ln |R_n|) \\ &\quad - \frac{1}{2} (\vec{y} - \vec{s}_m)' R_m (\vec{y} - \vec{s}_m) \\ &\quad + \frac{1}{2} (\vec{y} - \vec{s}_n)' R_n (\vec{y} - \vec{s}_n) \end{aligned} \quad (4-4)$$

and

$$\begin{aligned} I(w, a, m, n) &= \int \ln \left[\frac{f_m(\vec{y}|a)}{f_n(\vec{y}|a)} \right] f_m(\vec{y}|a) d\vec{y} \\ &= E \left\{ \ln \left[\frac{f_m(\vec{y}|a)}{f_n(\vec{y}|a)} \right] \right\} \end{aligned} \quad (4-5)$$

where the expectation is taken under the assumption $\vec{y} \sim N(\vec{s}_m, R_m)$.

Note that $I(w, a, m, m) = 0.0$ for $m = 1, \dots, L$.

Lemma 4.1 If $\vec{y} \sim N(\vec{c}, R)$ and R is positive definite, and A is symmetric, then

$$E\{\vec{y}' A \vec{y}\} = \text{tr}(AR^{-1}) + \vec{c}' A \vec{c}$$

Proof By theorem 10.3.2 of Graybill (1969)

$$E\{(\vec{y} - \vec{c})' A (\vec{y} - \vec{c})\}$$

$$\begin{aligned} &= \frac{|R|^{1/2}}{(2\pi)^{n/2}} \int_{-\infty}^{\infty} (\vec{y} - \vec{c})' A (\vec{y} - \vec{c}) e^{-1/2(\vec{y}-\vec{c})' R (\vec{y}-\vec{c})} d\vec{y} \\ &= \text{tr}(AR^{-1}) \end{aligned}$$

But

$$E\{(\vec{y} - \vec{c})' A (\vec{y} - \vec{c})\} = E\{\vec{y}' A \vec{y}\} - \vec{c}' A \vec{c}$$

The lemma follows immediately.

Q.E.D.

Applying the lemma to the expectations of the quadratic forms in equation (4-4) we see:

1. $\vec{y}_m \sim N(\vec{s}_m, R_m) \Rightarrow E\{(\vec{y}_m - \vec{s}_m)' R_m (\vec{y}_m - \vec{s}_m)\} = \text{tr}(R_m R_m^{-1}) = J$
2. $\vec{y} - \vec{s}_n \sim N(\vec{s}_m - \vec{s}_n, R_m) \Rightarrow E\{(\vec{y} - \vec{s}_n)' R_n (\vec{y} - \vec{s}_n)\}$

$$= \text{tr}(R_n R_m^{-1}) + (\vec{s}_m - \vec{s}_n)' R_n (\vec{s}_m - \vec{s}_n)$$

Thus

$$I(w, a, m, n) = \frac{1}{2} \left\{ \ln |R_m| - \ln |R_n| \right\} - \frac{1}{2} J + \frac{1}{2} \text{tr}(R_n R_m^{-1}) + \frac{1}{2} (\vec{s}_m - \vec{s}_n)' R_n (\vec{s}_m - \vec{s}_n) \quad (4-6)$$

$$\begin{aligned} I(w, a, m, n) + I(w, a, n, m) &= -J + \frac{1}{2} \left[\text{tr}(R_n R_m^{-1}) + \text{tr}(R_m R_n^{-1}) \right] \\ &+ \frac{1}{2} \left[(\vec{s}_m - \vec{s}_n)' R_n (\vec{s}_m - \vec{s}_n) + (\vec{s}_n - \vec{s}_m)' R_m (\vec{s}_n - \vec{s}_m) \right] \\ &= -J + \frac{1}{2} \left[\text{tr}(R_n R_m^{-1}) + \text{tr}(R_m R_n^{-1}) \right] \\ &+ \frac{1}{2} [(\vec{s}_m - \vec{s}_n)' (R_m + R_n) (\vec{s}_m - \vec{s}_n)] \end{aligned} \quad (4-7)$$

$$\begin{aligned} I(w, a) &= \sum_{n=2}^L \sum_{m=1}^{n-1} \theta_n \theta_m [I(w, a, m, n) + I(w, a, n, m)] \\ &= \sum_{n=2}^L \sum_{m=1}^{n-1} \theta_n \theta_m \left\{ -J + \frac{1}{2} \left[\text{tr}(R_n R_m^{-1}) + \text{tr}(R_m R_n^{-1}) \right] \right. \\ &\quad \left. + \frac{1}{2} [(\vec{s}_m - \vec{s}_n)' (R_m + R_n) (\vec{s}_m - \vec{s}_n)] \right\} \\ &= -J \sum_{n=2}^L \sum_{m=1}^{n-1} \theta_m \theta_n + \frac{1}{2} \sum_{n=1}^L \theta_n \text{tr} \left[\left(\sum_{m \neq n} \theta_m R_m \right) R_n^{-1} \right] \\ &\quad + \frac{1}{2} \sum_{n=2}^L \sum_{m=1}^{n-1} \theta_n \theta_m (\vec{s}_m - \vec{s}_n)' (R_m + R_n) (\vec{s}_m - \vec{s}_n) \end{aligned} \quad (4-8)$$

The last form of this equation appears to be the most convenient for computing purposes.

4.3 Intuitive Analysis

Looking at the computing form of equation (4-8) it can be seen that there are three terms. The first term is $-J \sum_{n=2}^L \sum_{m=1}^n \theta_m \theta_n$. The value of this term does not depend upon a and hence has no effect upon the choice of a . From this consideration we note that computing the value of this term would not be beneficial if only one more stage of experimentation is available.

The third term of the sum is a weighted sum of the quadratic forms

$$(\vec{s}_m - \vec{s}_n)' (R_m + R_n) (\vec{s}_m - \vec{s}_n)$$

Thus this term is in effect a separating function in the sense that these quadratic forms will be maximized when the pairs of expected values of \vec{y} under the various hypotheses are as far apart as possible in comparison to the precisions of \vec{y} . If the precisions R_m and R_n are large then \vec{s}_m and \vec{s}_n do not need to be far apart to provide much information whereas if these precisions are small then the expected values \vec{s}_m and \vec{s}_n must be further apart to provide the same information. The weighting factors are the products $\theta_n \theta_m$. Thus when θ_n and θ_m are both small, $\theta_n \theta_m$ is very small and the information due to the separation of \vec{s}_n and \vec{s}_m is discounted somewhat. If θ_n and θ_m are large then the information due to separation of \vec{s}_n and \vec{s}_m is given more importance. Thus this third term causes experiments to be chosen which separate the expected values of \vec{y} under the respective hypotheses which are still in serious contention for being chosen.

It is interesting to note that some authors (Hunt and Reiner (1965), e.g.) have proposed criteria for selection of experiments involving only distances between expected values. In a later paper, Box and Hill (1967) proposed that the distances as such are not important, but the distances weighted by some function of the variability about the expected values are important. It is seen here that the expected K-L information function does just that.

The second term in equation (4-8) is

$\frac{1}{2} \sum_{n=1}^L \theta_n \text{tr} \left[\left(\sum_{m \neq n} \theta_m R_m \right) R_n^{-1} \right]$. This can be thought of as a weighted sum of ratios of precisions. If only one y value is to be observed, this component becomes

$$\frac{1}{2} \sum_{n=1}^L \theta_n \frac{\sum_{m \neq n} \theta_m R_m}{R_n} \quad (4-9)$$

It would be interesting to see when this term is maximized. Upon taking partial derivatives of equation (4-9), setting to zero, and simplifying, one arrives at the following set of simultaneous non-linear equations.

$$\sum_{k=1}^L \theta_i \left(\frac{R_k^2 - R_i^2}{R_i} \right) = 0 \quad i = 1, \dots, L$$

It can be immediately seen that one solution to this system is

$R_1 = R_2 = \dots = R_L$. This solution implies that the experiments should tend to give the same precision for the expected value of \vec{y}

under each hypothesis. This term is not considered any further here.

In summary, it can be seen that the expected K-L information function in this case is basically a rather simple separating function. One would be hard pressed to construct a much simpler separating function which has more intuitive appeal. If multivariate observations are permitted, then it might be possible to delete the second term of equation (4-8) to save a good deal of computing.

CHAPTER 5

THE SEQUENTIAL DECISION PROCEDURE

Three components are required for a sequential adaptive decision procedure; (1) a rule which determines if sampling should be terminated or continued, (2) a rule which specifies the experiment to be performed given the current state of the system, and (3) a rule which selects the model equation which will be claimed to be true when sampling is terminated. The first part of this chapter discusses the experiment selection rule and the second section presents the stopping and model selection rules.

5.1 Experiment Selection Rule

The procedure adopted for this dissertation is the so-called myopic procedure. This rule simply chooses as the next experiment that one which maximizes the anticipated K-L information for the next stage only.

We assume that an upper limit, J_{MAX} , to the number of observations is specified. This number may be infinite. An allocation of the observations to the stages of sampling is described by a $J_{MAX} \times 1$ vector \vec{n} , where n_i gives the number of observations at stage i . The question arises as to how the observations should be allocated. That is, should all J_{MAX} be taken at once, strictly one-at-a-time, or in different sized groups. As the first step in answering this, let A_j denote the set of experiments in A which

specify that j observations should be taken. For any given state $w \in \Omega$, let $a_j^*(w)$ denote the element of A_j such that

$$I[w, a_j^*(w)] = \sup_{a_j \in A_j} I(w, a_j)$$

Lemma 5.1 For any $w \in \Omega$, and i, j such that $i > j$ we have

$$\underline{I[w, a_i^*(w)] \geq I[w, a_j^*(w)]}.$$

Proof: We introduce the following notation. Let $y_k(a_i^*)$, $k = 1, \dots, i$ denote the random variables observed under $a_i^*(w)$ and $y_k(a_j^*)$, $k = 1, \dots, j$ denote the random variables observed under a_j^* . Define another experiment $\tilde{a}_i \in A_i$ by choosing the first j observations according to a_j^* and the remaining $i - j$ observations according to the last $i - j$ of a_i^* . This leads to the random variables

$$\tilde{y}_k(\tilde{a}_i) = \begin{cases} y_k(a_j^*) & k = 1, \dots, j \\ y_k(a_i^*) & k = j + 1, \dots, i \end{cases}$$

Because $I(w, a, m, n)$ is positive definite and is additive for independent observations

$$I(w, \tilde{a}_i, m, n) \geq I(w, a_j^*, m, n)$$

Thus

$$\begin{aligned} I(w, \tilde{a}_i) &= \vec{\theta}' [I(w, \tilde{a}_i, m, n)] \vec{\theta} \\ &\geq \vec{\theta}' [I(w, a_j^*, m, n)] \vec{\theta} = I(w, a_j^*) \end{aligned}$$

But by definition $I(w, a_i^*) \geq I(w, \tilde{a}_i)$ and hence

$$I(w, a_i^*) \geq I(w, a_j^*)$$

Q.E.D.

The lemma simply proves that an experiment with more observations will be expected to provide more information than one with fewer observations. In determining an allocation one should also consider the cost of experimenting. In particular, if we assume that each observation has a constant cost associated with it, then it is reasonable to choose the experiment which maximizes

$$\frac{1}{j} I(w, a_j) \quad j = 1, \dots, J_{\text{MAX}}$$

Thus prior to stage k let $m = \sum_{i=1}^{k-1} n_i$ and assume $m < J_{\text{MAX}}$. The optimal experiment is the element $a^* \in A$ which for the current state w_{k-1} yields

$$j = 1, \dots, J_{\text{MAX}} - m \left\{ \max_{a \in A_j} \frac{1}{j} I(w_{k-1}, a) \right\}$$

If sampling has not been terminated by the rules developed in Chapter 5.2, then we stop when $\sum n_i = J_{\text{MAX}}$ and select the model according to the rules in Chapter 5.2.

5.2 Stopping and Model Selection Rules

We now discuss the problems of determining which of the postulated models is the true one and determining when the results of the experiments are sufficiently informative to stop sampling and make the choice.

Box and Hill (1967) suggested that for their procedure, experimenting be terminated whenever one model is clearly superior to

the others. This is obviously a reasonable statement but it is in need of formal definition before it can be used as a stopping and selection rule. We propose general stopping and selection rules and a modified version which might be used in certain instances involving nested models.

(1) Stopping rule: Let θ_m be some specified value $1/L < \theta_m \leq 1.0$. Let J_{MAX} denote the maximum number of observations permitted. Then terminate sampling whenever either $\max_{i=1,L} \{\theta_i\} \geq \theta_m$ or J_{MAX} observations have been taken, whichever occurs first.

(2) Model selection rule: Upon termination choose the correct model to be H_j^* where $\theta_j^* = \max_{i=1,L} \{\theta_i\}$.

We now present a modified stopping and selection procedure for use with nested models which may be of some value when θ_m is very near 1.0 and/or when J_{MAX} is relatively large. The reason for presenting a modified procedure arises from the large sample results of Chapter 3 and the Monte-Carlo results of Chapter 6. First, if H_{i^*} denotes the unknown true model, it is not known whether $\theta_{i^*,j} \rightarrow 1.0$ or not. From the Monte-Carlo results it seems that the typical behavior of $\theta_{i^*,j}$ for nested models is to fairly rapidly increase to something near 1.0 and then fluctuate, possibly slowly approaching 1.0. Thus, if θ_m is very near 1.0 it may be that extremely large samples would be required. Thus we would like to reduce the average sample size without seriously detracting from the probability of choosing the correct model.

To introduce the modified procedure consider the following example:

$$H_1: y = \beta_1 x_1 + \varepsilon$$

$$H_2: y = \beta_1 x_1 + \beta_2 x_2 + \varepsilon$$

If H_1 is actually the true model, then the posterior distribution of (β_1, β_2) under H_2 should approach a point distribution with $\mu_2^{(2)} = 0$ and $\mu_1^{(2)}$ equal to the unknown value of the parameter. However, $\theta_{1,j}$ may not approach 1.0. Assume some small positive constant γ is specified. Then after each stage of sampling, test if

$$d = \left[\mu_1^{(1)} - \mu_1^{(2)} \right]^2 + \left[\mu_2^{(2)} \right]^2 \leq \gamma$$

If $d \leq \gamma$ then drop model 2 from contention and replace θ_1 by $\theta_1 + \theta_2$. Then apply the previously described stopping and selection rules. In this simple example, the dropping of model 2 would automatically cause sampling to be terminated. This would of course not necessarily be true in more general situations.

To generalize the procedure some additional notation and concepts must be introduced. We use the symbol \supset to denote inclusion. Thus $H_i \supset H_j$ means that the model of H_j is nested within the model of H_i . The set of models $\{H_i\}$ is a partially ordered set under the partial ordering relation \supset . In the theory of partially ordered sets a chain is defined to be a partially ordered set such that for any two elements (H_1 and H_2 say) of the set either $H_1 \supset H_2$ or $H_2 \supset H_1$. For the purposes of this dissertation we define a string of elements from the partially ordered set as a sub-

set of elements such that the subset forms a chain. A maximal string is constructed from any string by adding all the elements of $\{H_i\}$ to the string which can be added without causing the enlarged set to lose the property of being a chain.

To formulate the modified stopping and selection procedure we first construct all of the maximal strings that can be constructed from the set $\{H_i\}$ and order the elements of the strings using the relation \supset .

For example, suppose $L = 5$ and the five models are as specified below:

Model number	Model equation
1	$y = \beta_0 + \beta_1 x_1 + \epsilon$
2	$y = \beta_0 + \beta_3 x_3 + \epsilon$
3	$y = \beta_0 + \beta_1 x_1 + \beta_2 x_2 + \epsilon$
4	$y = \beta_0 + \beta_1 x_1 + \beta_3 x_3 + \epsilon$
5	$y = \beta_0 + \beta_1 x_1 + \beta_2 x_2 + \beta_3 x_3 + \epsilon$

The maximal strings are easily verified as being

$$(H_1 \subset H_3 \subset H_5)$$

$$(H_1 \subset H_4 \subset H_5)$$

$$(H_2 \subset H_4 \subset H_5)$$

In each of these strings the maximal element is H_5 .

The modified procedure consists of computing for the first

maximal string the squared distance of the posterior mean vector of the maximal element from the posterior mean vector of the submaximal or next largest submodel of the string. If this quantity is less than some prespecified value γ , the maximal model is dropped from the set $\{H_i\}$. The posterior probability of the maximal element is added to the probability of the next element of the string until either only one model remains or there is no need to drop models. Before considering the next maximal string, all models which have been dropped must also be deleted from the remaining strings. The above procedure is then repeated for each maximal string in turn.

Once this has been completed and all models which can be dropped because they reduce to models with fewer parameters have been dropped, the same stopping and selection rules proposed for the non-nested case are applied.

Note: The procedure just described is not necessarily the best or the most natural one to use for combining models. For example, an alternative to the distance of the means might be to combine models when the probability distribution of the maximal element is sufficiently concentrated about the mean of the submaximal element. This would have the advantage of using the information contained in the precisions of the distributions also.

CHAPTER 6

COMPUTER SIMULATION RESULTS

The purpose of this chapter is to report and discuss the results of a number of Monte-Carlo simulation studies of the sequential procedure proposed in Chapter 5. The chapter is divided into four major sections. The first section describes the general simulation procedure and presents a brief description of the algorithm used. A computer program based on this algorithm is described in further detail in appendix A. The second section presents and discusses the results of a number of simulations performed to gain further information about the large sample behavior first discussed in Chapter 3. The primary concern is with the posterior probabilities and the posterior means of the parameter distributions after a large number of observations. The third section presents and discusses the results of some simulation studies of the proposed sequential procedure when the number of observations is limited and when the stopping rule of Chapter 5 is used. The primary concern is with the probability of the procedure actually selecting the correct model and the average sample size required until termination. The last section of the chapter presents an overall discussion of the results.

6.1 General Simulation Procedure

The sequential procedure proposed in Chapter 5 consisted of

(1) an experiment termination rule, (2) an experiment selection rule, and (3) a model selection rule. Because of the mathematical complexity of the posterior distributions involved it was not feasible to analytically examine how well these rules work. The general procedure by which the Monte Carlo simulation technique was used to study performance is outlined in the following algorithm.

1. Input:

$\vec{\mu}_{k,0}$ the prior means of the parameters of the models
 $\Psi_{k,0}$ the prior precision matrices of the parameters
 of the models
 $\theta_{k,0}$ the prior probabilities of the models being
 correct
 N the number of simulations
 θ_m stopping probability
 J_{MAX} maximum number of observations
 i^* the model chosen to generate the observed variable
 $\vec{\mu}^*$ values of the parameters of the true model

2. $n \leftarrow 0$

3. $PCS \leftarrow 0$

4. $\bar{N}_i \leftarrow 0$ (for $i = 1, J_{MAX}$)

5. $j \leftarrow 0$

6. $j \leftarrow j + 1$

7. Determine optimal a^* as described in Chapter 4. Denote as a^* and let M_{a^*} denote design matrix for model i^* when a^* is chosen. (All simulations in this dissertation

consider strictly one-at-a-time sampling for simplicity.)

8. $y_j \leftarrow M_{a^*}^{\rightarrow*}$
9. Generate a pseudo-random observation ϵ_j from a $N(0, \tau)$ distribution. (Described in detail in appendix A)
10. $y_j \leftarrow y_j + \epsilon_j$
11. For $\ell = 1, \dots, L$ compute $\theta_{\ell,j}$, $\psi_{\ell,j}$, and $\vec{\mu}_{\ell,j}$ from y_j and $\theta_{\ell,j-1}$, $\psi_{\ell,j-1}$, and $\vec{\mu}_{\ell,j-1}$ as described in Chapter 3.
12. Find k such that $\theta_{k,j} = \text{MAX}_i \{\theta_{i,j}\}$
13. If $j \geq J_{\text{MAX}}$ or $\theta_{k,j} \geq \theta_m$ go to 14. Otherwise go to 6.
14. $\bar{N}_j \leftarrow \bar{N}_j + 1$
15. If $k = i^*$; $\text{PCS} \leftarrow \text{PCS} + 1$
16. $n \leftarrow n + 1$
17. If $n \geq N$ go to 18. Otherwise go to 5.
18. $\text{PCS} \leftarrow \text{PCS}/N$
19. $\text{ASN} \leftarrow \left(\sum_{i=1}^{J_{\text{MAX}}} i \bar{N}_i \right) / N$
20. Stop

Upon stopping, the value of PCS is the observed probability of correctly choosing i^* as the true model for the prior distributions specified when in fact the true value of the parameters is given by $\vec{\mu}^*$. ASN gives the average sample number upon termination.

The above algorithm can be easily used for either large sample or small sample studies. For example, for large sample studies set $\theta_m = 1.0$, $N = 1$, and J_{MAX} to some large number, say 100 or 500.

For small sample studies set $\theta_m < 1.0$, J_{MAX} to some small number, and N to some larger number, say 500 or 1000.

6.2 Large Sample Studies

In this section we examine the large sample properties of the posterior probabilities of the models and the posterior means of the parameter distributions. Three sets of problems are studied. First, two sets of nested polynomial models are studied. The posterior probabilities of each model, the posterior means of the parameter distributions, and the proportion of times each of the allowable values of the independent variable is chosen as optimal are tabulated for simulations of 100 and 500 observations. Second, one set of nested factorial models is studied for three different prior distributions on the models. And third, one set of non-nested factorial models is simulated. For the last two, the posterior probabilities and means of the parameter distributions are tabulated.

6.2.1 Polynomial Model Studies

Two sets of nested polynomial models are considered which have the following general form:

$$H_\ell: y = \sum_{j=0}^{\ell-1} \beta_j x^j + \epsilon, \ell = 1, L$$

Two values of L are studied, and for each of these choices, two choices of H_{1*} are made. The values of τ , $\theta_{\ell,0}$, and $\psi_{\ell,0}$ are specified as

$$\tau = 100.0$$

$$\psi_{\ell,0} = 1$$

$$\theta_{\ell,0} = \frac{1}{L}$$

for all simulations. The values of $\vec{\mu}_{\ell,0}$ are tabulated at the tops of figures 1 and 2 and the resulting functions are graphed on the interval $x \in [-1,+1]$ at the bottoms of the respective figures. For $L = 4$, the two choices of H_{i^*} are H_2 and H_3 . For $L = 6$, the two choices of H_{i^*} are H_3 and H_5 . For simplicity, the actual values of the parameters used to generate the data were chosen to be $\vec{\mu}_{i^*,0}$ for each of the four cases.

For these simulations, the definition of A was arbitrarily taken to be

$$A = \{a^{(i)}: i = 0, \dots, 9\}$$

where

$$a^{(i)}: x = -1 + \frac{2i}{9}$$

Note that sampling is strictly one observation per stage.

The simulation results are summarized in table 1 and given in further detail in tables 2 through 9. For each choice of L and i^* , five simulations of 100 observations and five simulations of 500 observations were performed. For these simulations, the sample paths of the $\theta_{\ell,j}$ were printed out and the choice of $a^{(i)}$ at each stage were printed. The posterior means of the parameter distributions were printed only after the last stage. Tables 2, 4, 6, and 8 give the posterior probabilities after 100 observations and the first 100 out of 500 observations. The proportions p_i of using $a^{(i)}$ are also given. Tables 3, 5, 7, and 9 give the same informa-

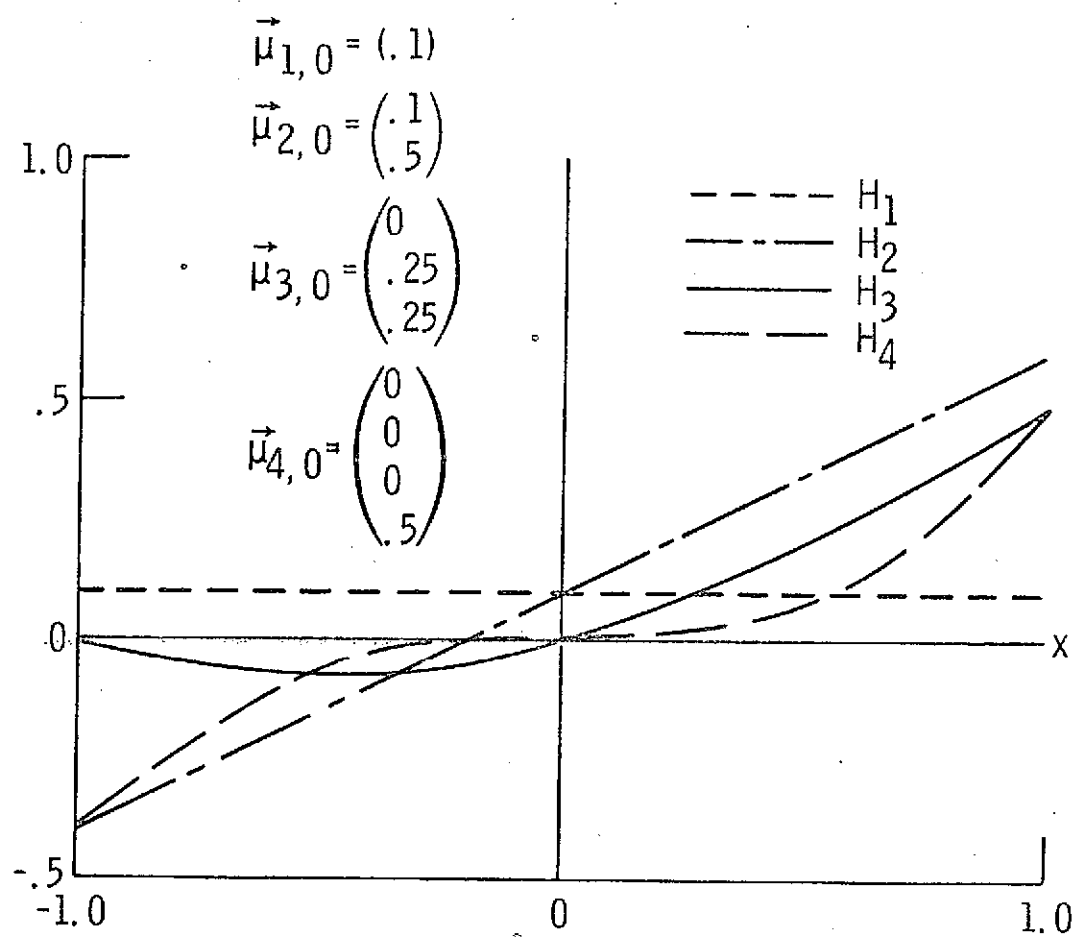


Figure 1. - Tabulations of the prior means of the parameters and graphs of the resulting functions over the interval $[-1, +1]$ for large sample polynomial study one.

$$\begin{aligned}
 \vec{\mu}_{1,0} &= (0) \\
 \vec{\mu}_{2,0} &= \begin{pmatrix} 0 \\ .5 \end{pmatrix} \\
 \vec{\mu}_{3,0} &= \begin{pmatrix} 0 \\ .25 \\ .25 \end{pmatrix} \\
 \vec{\mu}_{4,0} &= \begin{pmatrix} 0 \\ 0 \\ 0 \\ .5 \end{pmatrix} \\
 \vec{\mu}_{5,0} &= \begin{pmatrix} 0 \\ .5 \\ 0 \\ 0 \\ .1 \end{pmatrix} \\
 \vec{\mu}_{6,0} &= \begin{pmatrix} 0 \\ .5 \\ 0 \\ 0 \\ 0 \\ .1 \end{pmatrix}
 \end{aligned}$$

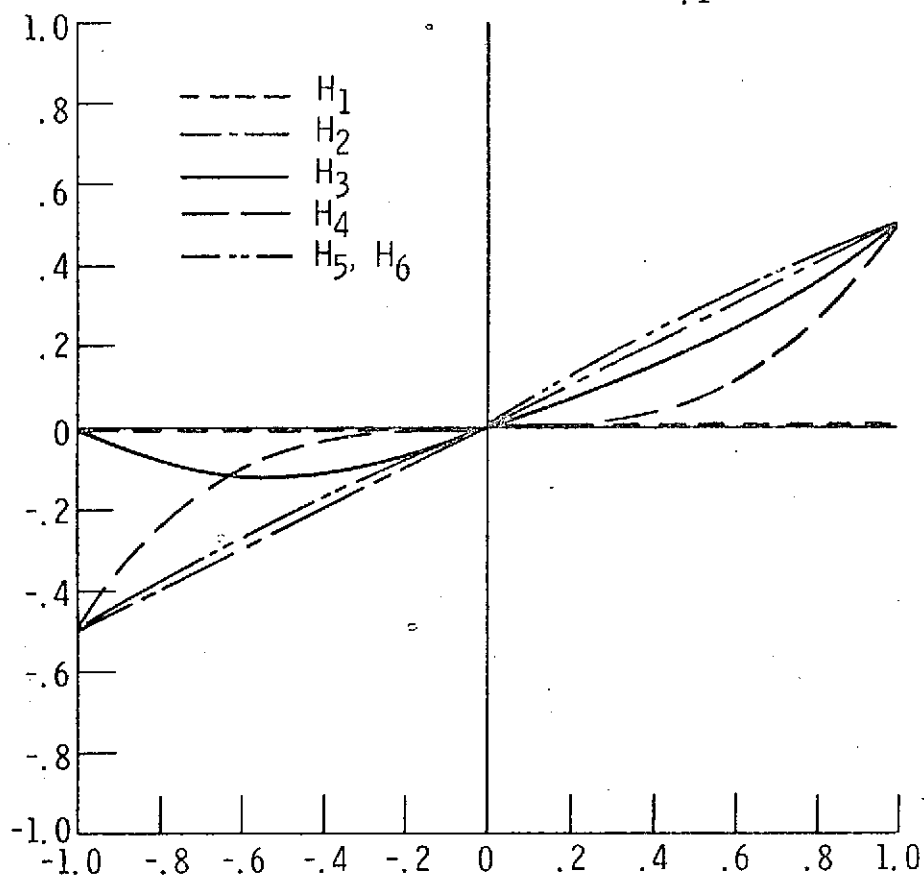


Figure 2. - Tabulations of the prior means of the parameters and graphs of the resulting functions over the interval $[-1, +1]$ for large sample polynomial study two.

tion for the 500 observation simulations.

Figures 3 and 4 present typical sample paths for the posterior probability of the correct model. In figure 3, the value of $\theta_{2,j}$ is plotted for the first 250 observations of the third simulation for $L = 4$ and $i^* = 2$. In figure 4, the value of $\theta_{3,j}$ is plotted for the first 250 observations of the first simulation for $L = 4$ and $i^* = 3$. These figures illustrate the typical behavior of $\theta_{i^*,j}$. It fairly rapidly rises to a value of about 0.85 to 0.95 and then slowly and erratically oscillates. This is suspected to be because of the nested nature of the model equations. It was because of this behavior that the modified selection rule of Chapter 5 was first introduced. Consideration of the posterior means of the parameter distributions will also provide some information concerning this modified rule.

For $L = 4$, consideration of tables 2, 3, 4, and 5 show that as j increases, $\vec{\mu}_{i,j} \rightarrow \begin{pmatrix} \mu_{i^*,j} \\ 0 \end{pmatrix}$ for $i > i^*$. This is in accord with the conclusions of Chapter 3. For $L = 6$ and $i^* = 3$ we again see the same close agreement with Chapter 3 as evidenced by tables 6 and 7. However, for $i^* = 5$, an entirely different situation arises. To understand this we should note that the model used to generate the sequential observations is

$$y = 0.5x + 0.1x^4 + \varepsilon$$

This function can be very closely approximated by a model of the form

$$y = ax + bx^2 + \varepsilon$$

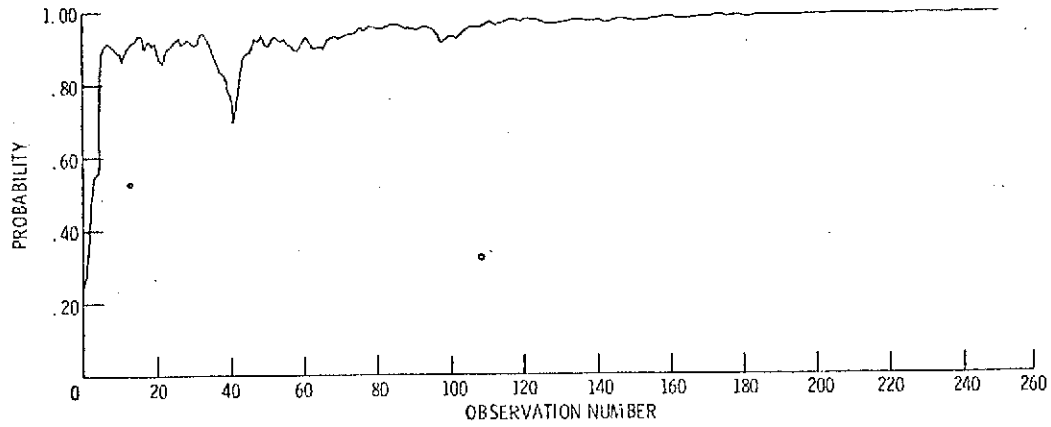


Figure 3. - The sample path of $\theta_{2,j}$ for $L = 4$, $i^* = 2$ for the first 250 observations of simulation no. 3. A well behaved path for nested models.

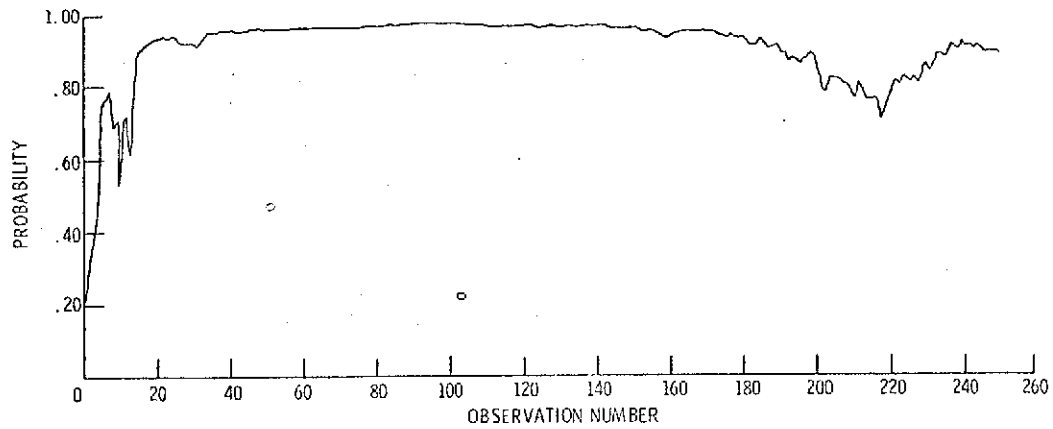


Figure 4. - The sample path of $B_{3,j}$ for $L = 4$, $i^* = 3$ for the first 250 observations of simulation no. 1. A typical path for nested models.

over the range of x values considered. And in fact we note that there is a marked preference for choosing the lower degree model as indicated by $\theta_{3,j}$ becoming close to 1.0. It is also interesting to note the behavior of $\vec{\mu}_{i,j}$ for $i > 3$. We do not see that $\vec{\mu}_{1,j} \rightarrow \begin{pmatrix} \vec{\mu}_{3,j} \\ 0 \end{pmatrix}$ as might be expected when H_3 is so close to being true, except for the case of $i = 4$. For $\vec{\mu}_5$ we note that the average posterior mean of the coefficient of x^3 is quite close to zero and the sum of the posterior means of the coefficients of x^2 and x^4 is quite close to 0.1. For $\vec{\mu}_6$ we note that the sums of the posterior means of the coefficients of x^2 and x^4 is close to 0.1 and the sum of the posterior means of the coefficients of x , x^3 , and x^5 is close to 0.5. From these simulation studies it is not clear whether this behavior is simply because 500 observations is not a sufficiently large number to discriminate well between such nearly equivalent functions or if this behavior will persist no matter how large the number of observations.

We now turn to a discussion of the observed proportions of times the $a^{(i)}$ were chosen as the optimal experiments. From tables 2 and 3 which present the results of $L = 4$ and $i^* = 2$ we see that the largest p_i are for p_0 , p_4 , p_5 , and p_9 . These correspond to $x = -1$, $x = -1/9$, $x = +1/9$, and $x = +1$. Because of the discretization of the interval $(-1, +1)$ we might assume that the asymptotically most informative experiments were $x = -1$, $x = 0$, and $x = +1$. From tables 4 and 5 we see the largest p_i are p_0 , p_2 , p_7 , and p_9 corresponding to $x = -1$, $x = -5/9$, $x = +5/9$, and

$x = +1$. The relationship of these proportions and x points to the experimental designs which are optimal from other considerations might be interesting. For example, Kiefer and Wolfowitz (1959) consider optimal designs for regression problems of a somewhat different nature. The comparison of the current results with such other works is currently being pursued but will not be reported in this dissertation.

6.2.2 Nested Factorial Models

A second set of simulation studies were made using the following models

$$H_1: y = \beta_0 + \epsilon$$

$$H_2: y = \beta_0 + \beta_1 x_1 + \epsilon$$

$$H_3: y = \beta_0 + \beta_1 x_1 + \beta_2 x_2 + \epsilon$$

$$H_4: y = \beta_0 + \beta_1 x_1 + \beta_2 x_2 + \beta_3 x_1 x_2 + \epsilon$$

with $\tau = 2.0$, $\psi_{\ell,0} = 1$, and $i^* = 3$. The prior means $\vec{\mu}_{\ell,0}$ were chosen as

$$\vec{\mu}_{1,0}^* = (0)$$

$$\vec{\mu}_{2,0}^* = (0,1)$$

$$\vec{\mu}_{3,0}^* = (0,1,-1) = (\vec{\mu}^*)^b$$

$$\vec{\mu}_{4,0}^* = (0,1,-1,0)$$

Three sets of $\theta_{\ell,0}$ were chosen:

$$1. \quad \theta_{1,0} = 0.1$$

$$\theta_{2,0} = 0.2$$

$$\theta_{3,0} = 0.3$$

$$\theta_{4,0} = 0.4$$

$$2. \quad \theta_{\ell,0} = 0.25 \quad \ell = 1,4$$

$$3. \quad \theta_{1,0} = 0.4$$

$$\theta_{2,0} = 0.3$$

$$\theta_{3,0} = 0.2$$

$$\theta_{4,0} = 0.1$$

The experiment space A is defined as $A = \{(x_1, x_2): x_i = \pm 1\}$. Note that experimenting is strictly one-at-a-time. The sequential selection procedure of Chapter 5 was used for five simulations of 500 observations each. The results are presented in table 10. We note that the posterior values are again in close agreement with Chapter 3 and the results of the polynomial models. There does not seem to be a pronounced effect upon the posterior probabilities of the models from changing the prior distribution although there does appear to be slightly higher posterior values of $\theta_{3,500}$ when the $\theta_{\ell,0}$ distribution is skewed toward the lower values. A possible explanation for this is that when this distribution is skewed toward the high values, the procedure is choosing experiments primarily to discriminate between H_3 and H_4 . Since H_3 is true, the model of

H_4 will rapidly become close to that of H_3 and the resulting experiments will not be very informative. When the prior probabilities are larger for the lower degree polynomials, however, the procedure chooses experiments primarily to discriminate between H_1 and H_2 . These experiments should then more rapidly tend to prove H_1 and H_2 to be inadequate.

6.2.3 Non-Nested Factorial Models

In this study, the following non-nested models were studied.

$$H_1: y = \beta_0 + \beta_1 x_1 + \beta_2 x_2 + \epsilon$$

$$H_2: y = \beta_0 + \beta_1 x_1 + \beta_3 x_3 + \epsilon$$

$$H_3: y = \beta_0 + \beta_2 x_2 + \beta_3 x_3 + \epsilon$$

The values of the parameters are chosen as

$$\left. \begin{aligned} i^* &= 1 \\ \vec{\mu}_{\ell,0} &= (1,1,1)^T \\ \psi_{\ell,0} &= 1 \\ \theta_{\ell,0} &= \frac{1}{3} \end{aligned} \right\} \ell = 1, \dots, L$$

$$\tau = 0.01, 1, 100$$

$$\vec{\mu}^* = (1,1,1)^T$$

The experiment space A was assumed to allow only one observation at a time with $x_i = \pm 1$. Five simulations were performed for each value of τ . For $\tau = 100$ it took only three observations for θ_3 to become 1.0 (within the accuracy of the computer). For $\tau = 1.0$ the number of observations required for the final posterior probabilities to reach 1.0 are tabulated in table 11. For $\tau = 0.01$,

1000 observations were taken and the resulting posterior probabilities are given in table 11. Again the results are in general agreement with Chapter 3.

6.3 Small Sample Performance Studies

In this section we examine the performance of the proposed sequential procedure as measured by the PCS and ASN values. First, two studies are presented of the problem of discriminating among the three models

$$H_1: y = \beta_1 x_1 + \epsilon$$

$$H_2: y = \beta_2 x_2 + \epsilon$$

$$H_3: y = \beta_1 x_1 + \beta_2 x_2 + \epsilon$$

The first study assumes H_3 is true and the second study assumes H_2 is true. The experiment space A is defined as

$$A = \{(x_1, x_2): x_i = \pm 1; \text{one-at-a-time sampling}\}$$

Then we consider the problem of choosing among the four nested models.

$$H_1: y = \beta_0 + \epsilon$$

$$H_2: y = \beta_0 + \beta_1 x_1 + \epsilon$$

$$H_3: y = \beta_0 + \beta_1 x_1 + \beta_2 x_2 + \epsilon$$

$$H_4: y = \beta_0 + \beta_1 x_1 + \beta_2 x_2 + \beta_3 x_1 x_2 + \epsilon$$

where A is as in the first two studies.

The primary reasons for concentrating on these simple models are that the parameter spaces are of a low enough dimension that

they can be visualized and they are small enough that extensive simulation studies would not require inordinate amounts of computer time. The fact that the parameter spaces can be visualized allows the effect of varying prior means upon PCS and ASN to be more easily grasped. Note that the modified stopping and selection rule was not used.

6.3.1 Study One - H_3 Assumed True

We study discriminating among

$$H_1: y = \beta_1 x_1 + \epsilon$$

$$H_2: y = \beta_2 x_2 + \epsilon$$

$$H_3: y = \beta_1 x_1 + \beta_2 x_2 + \epsilon$$

$$A = \{(x_1, x_2): x_1 = \pm 1; \text{one-at-a-time sampling}\}$$

where

$$\psi_{\ell,0} = 1 \quad \theta_{\ell,0} = \frac{1}{3}$$

$$\vec{\mu}_{1,0} = (1.0) \quad \vec{\mu}_{2,0} = (1.0)$$

and

$$\vec{\mu}^* = \begin{pmatrix} 1.0 \\ 1.0 \end{pmatrix}$$

Then a number of simulation experiments were performed for each combination of:

$$\tau = 0.50, 1.0, 2.0$$

$$\theta_m = 0.70, 0.80, 0.90$$

$$J_{\text{MAX}} = 8, 16$$

$$\vec{\mu}_{3,0} = \begin{pmatrix} 0.0 \\ 0.0 \end{pmatrix}, \begin{pmatrix} 0.50 \\ 0.50 \end{pmatrix}, \begin{pmatrix} 1.0 \\ 1.0 \end{pmatrix}, \begin{pmatrix} 1.5 \\ 1.5 \end{pmatrix}$$

The experiments for $J_{\text{MAX}} = 8$ used 1500 simulations and for $J_{\text{MAX}} = 16$ used 1000 simulations.

The choice of prior means deserves some comment. Figure 5 illustrates the points in (β_1, β_2) coordinate space corresponding to the prior means. The points corresponding to $\vec{\mu}_{1,0}$ and $\vec{\mu}_{2,0}$ are as close to $\vec{\mu}^*$ as possible since $\vec{\mu}_{1,0}$ is restricted to the horizontal axis and $\vec{\mu}_{2,0}$ to the vertical. The four choices for $\vec{\mu}_{3,0}$ then span a range about $\vec{\mu}^*$ and hence the resulting PCS and ASN values will indicate the importance of mis-specified prior means.

Tables 12 and 13 present the observed PCS and ASN values for the combinations of θ_m , τ , and $\vec{\mu}_{3,0}$. These results are also plotted as parametric surfaces in figures 6 through 9.

In general, the results are about what should be expected. The PCS increases with τ and ASN decreases with τ . PCS increases as $\vec{\mu}_{3,0}$ gets closer to $\vec{\mu}^*$. We also note that in most cases, PCS increases with θ_m for fixed τ and $\vec{\mu}_{3,0}$. There is, however, a distinct dropping off of PCS with θ_m along the peaks of the surfaces. There does not seem to be any ready explanation for this.

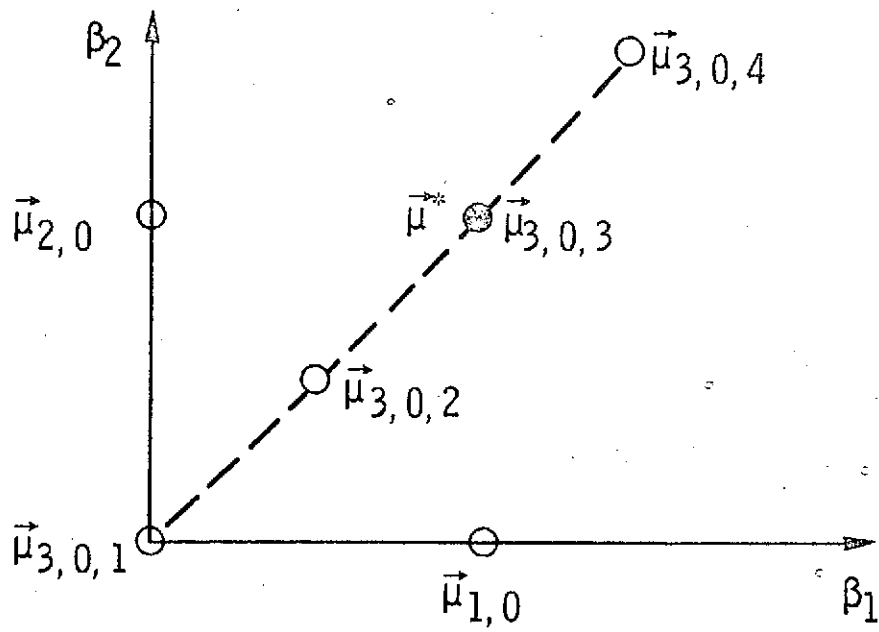


Figure 5. - Illustration of prior means for performance simulation experiment one.

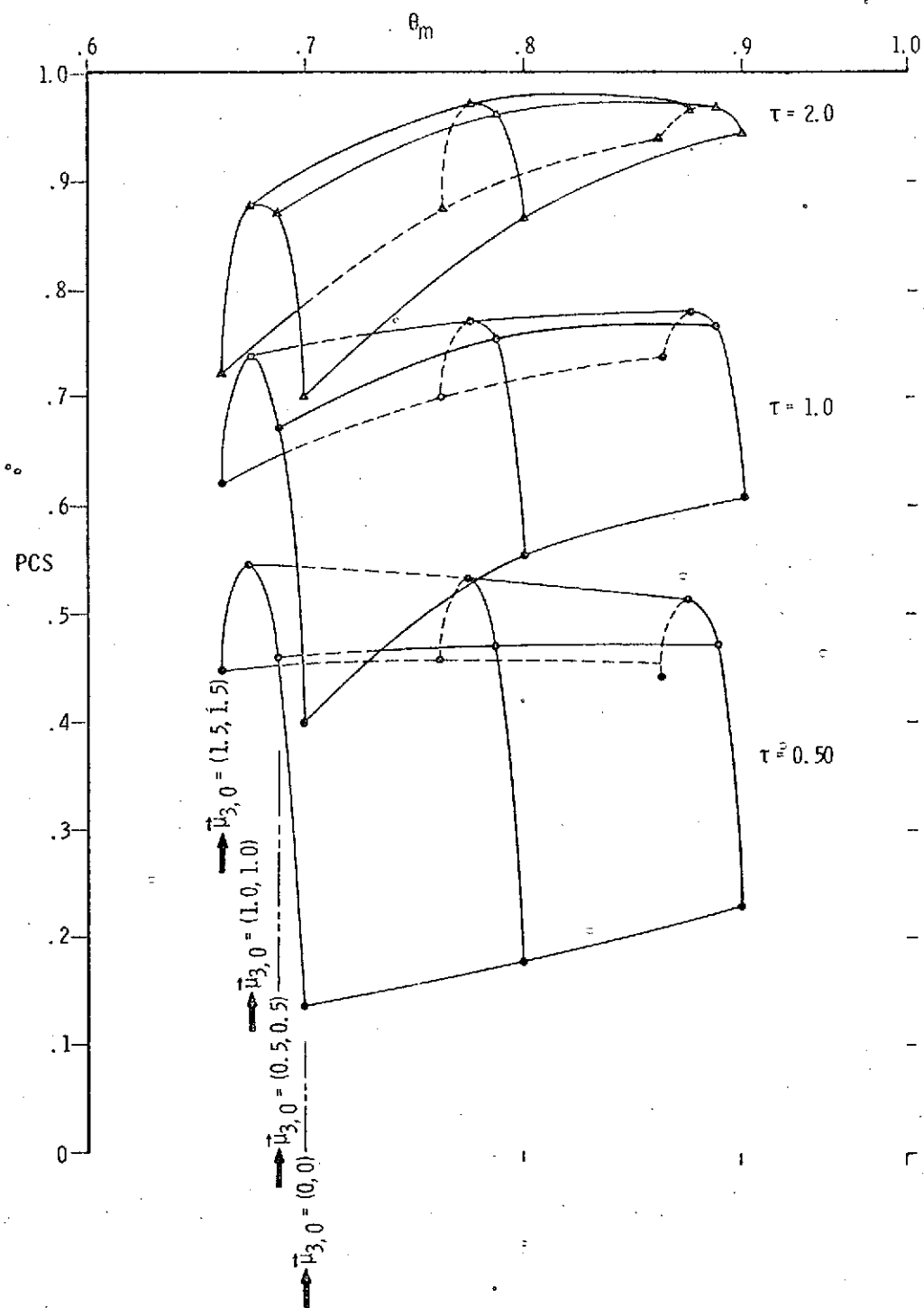


Figure 6. - Probability of correct selection (PCS) as a function of θ_m , τ , $\vec{\mu}_{3,0}$ for $J_{\max} = 8$ and H_3 true. Small sample performance simulation experiment one.

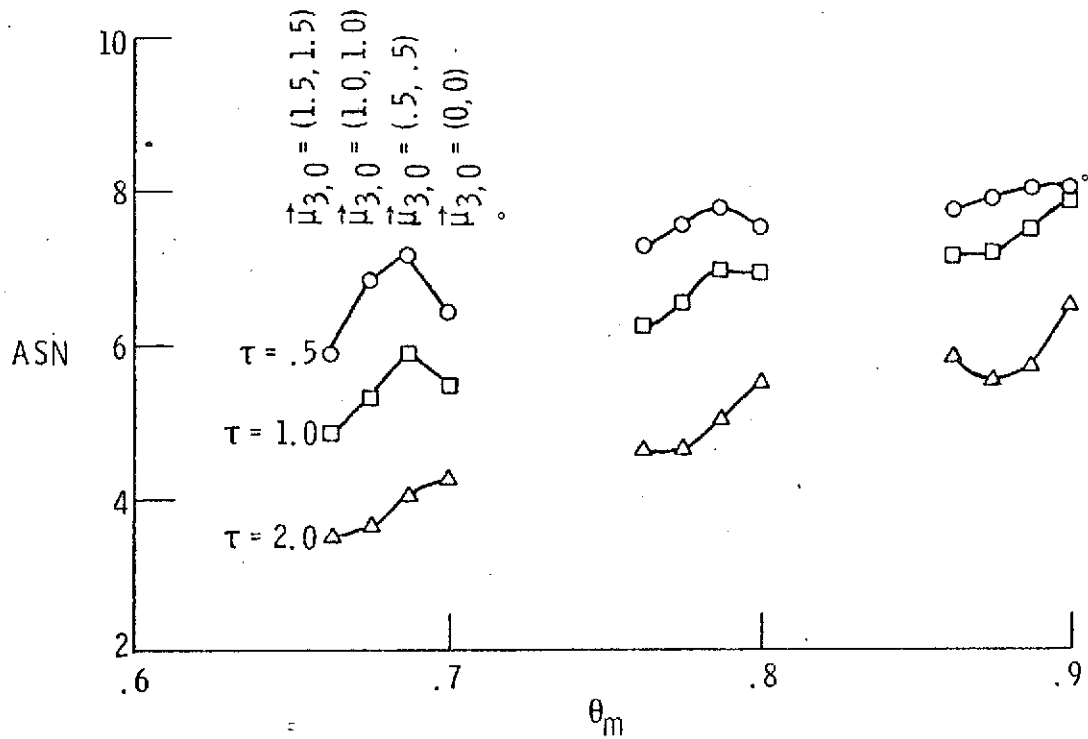


Figure 7. - Average sample number (ASN) as a function of θ_m , τ , $\vec{\mu}_{3,0}$ for $J_{\max} = 8$ and H_3 true. Small sample performance simulation experiment one.

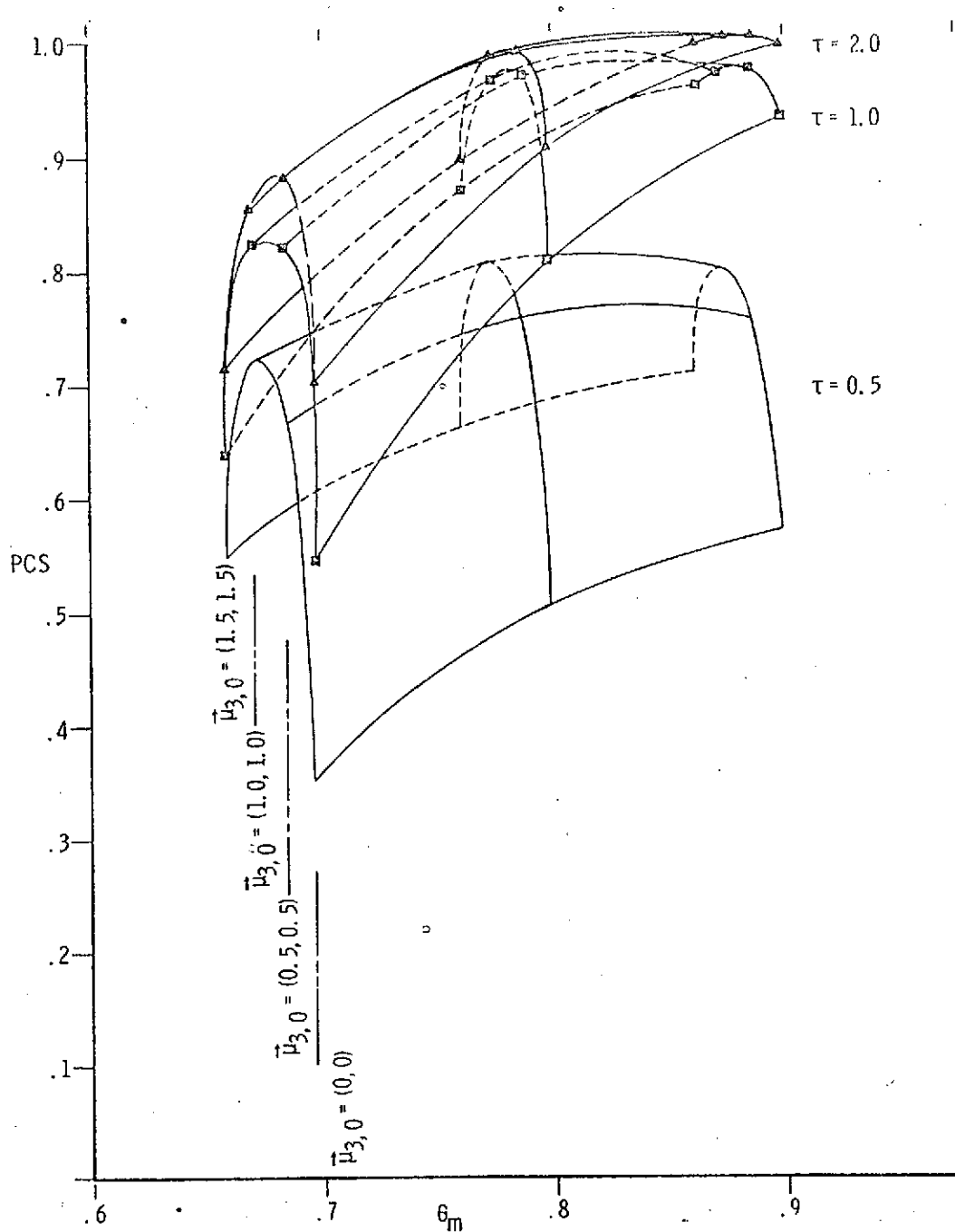


Figure 8. - Probability of correct selection (PCS) as a function of θ_m , τ , $\mu_{3,0}$ for $J_{\max} = 16$ and H_3 true. Small sample performance simulation experiment one.

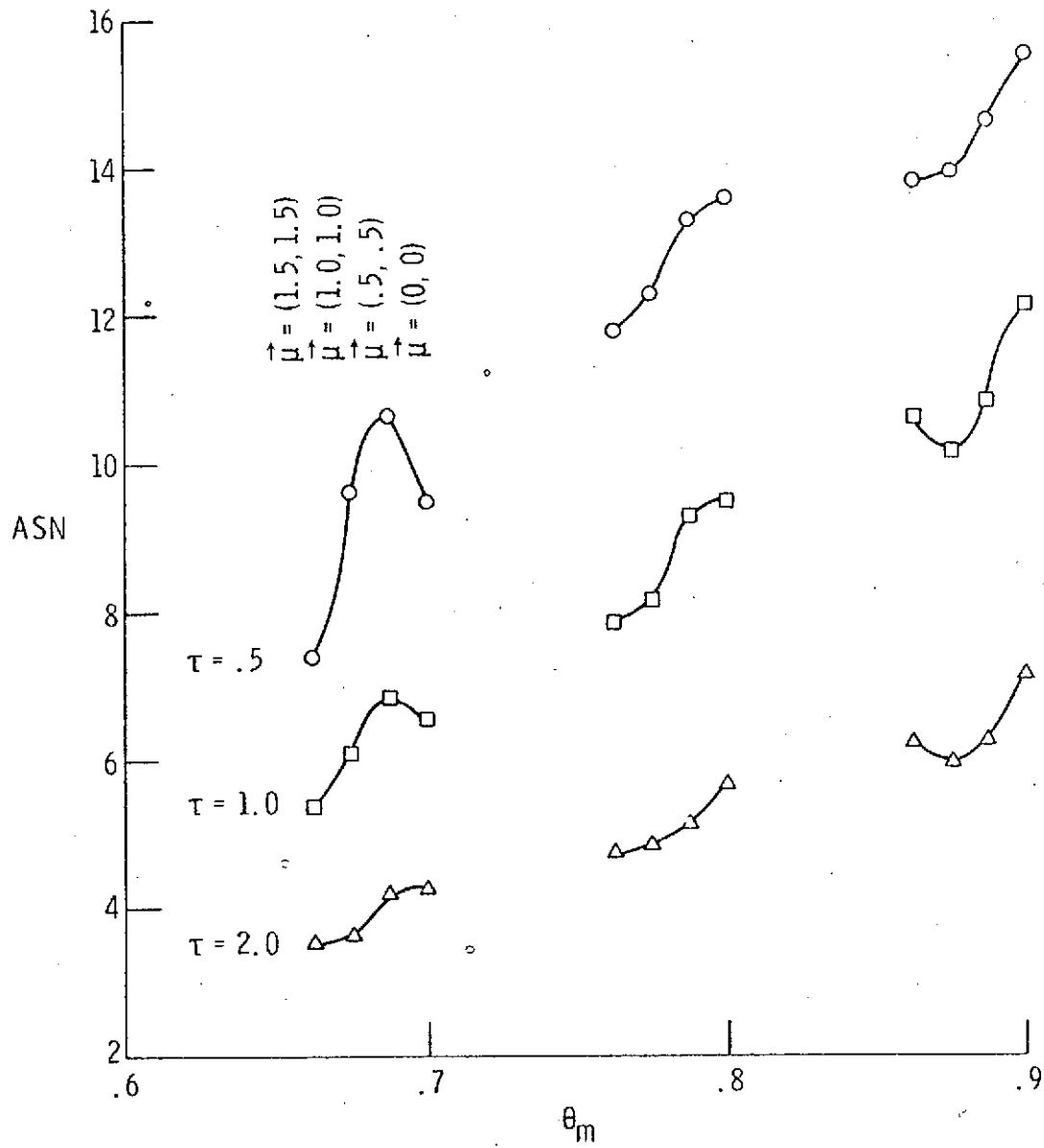


Figure 9. - Average sample number (ASN) as a function of θ_m , τ , $\bar{\mu}_{3,0}$ for $J_{\max} = 16$ and H_3 true. Small sample performance simulation experiment one.

6.3.2 H_2 Assumed True

A much less extensive study of this case was made than the case of H_3 assumed true. The same model equations were postulated and we assume

$$\left. \begin{aligned} \psi_{\ell,0} &= 1 \\ \theta_{\ell,0} &= \frac{1}{3} \end{aligned} \right\} \quad \ell = 1, \dots, L$$

$$\vec{\mu}_{1,0} = (0.0)$$

$$\vec{\mu}_{3,0} = (0.0, 1.0)'$$

$$\vec{\mu}^* = (1.0)$$

The values of τ , θ_m , and $\vec{\mu}_{2,0}$ which were simulated are tabulated in table 14 along with the simulation results. Figure 10 illustrates the prior means. Only one level of J_{MAX} (=8) was considered. Also, only 500 simulations were performed for each of these cases. The PCS results are also graphed as a parametric surface in figure 11. The results are generally the same as for H_3 true.

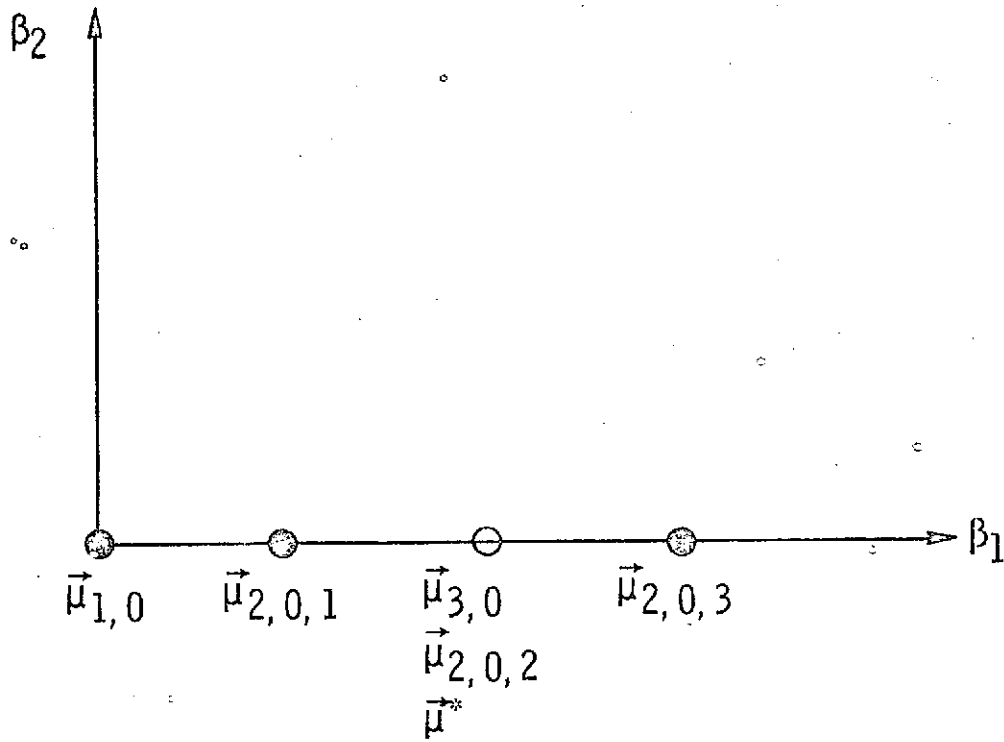


Figure 10. - Illustration of prior means for small sample performance simulation experiment two.

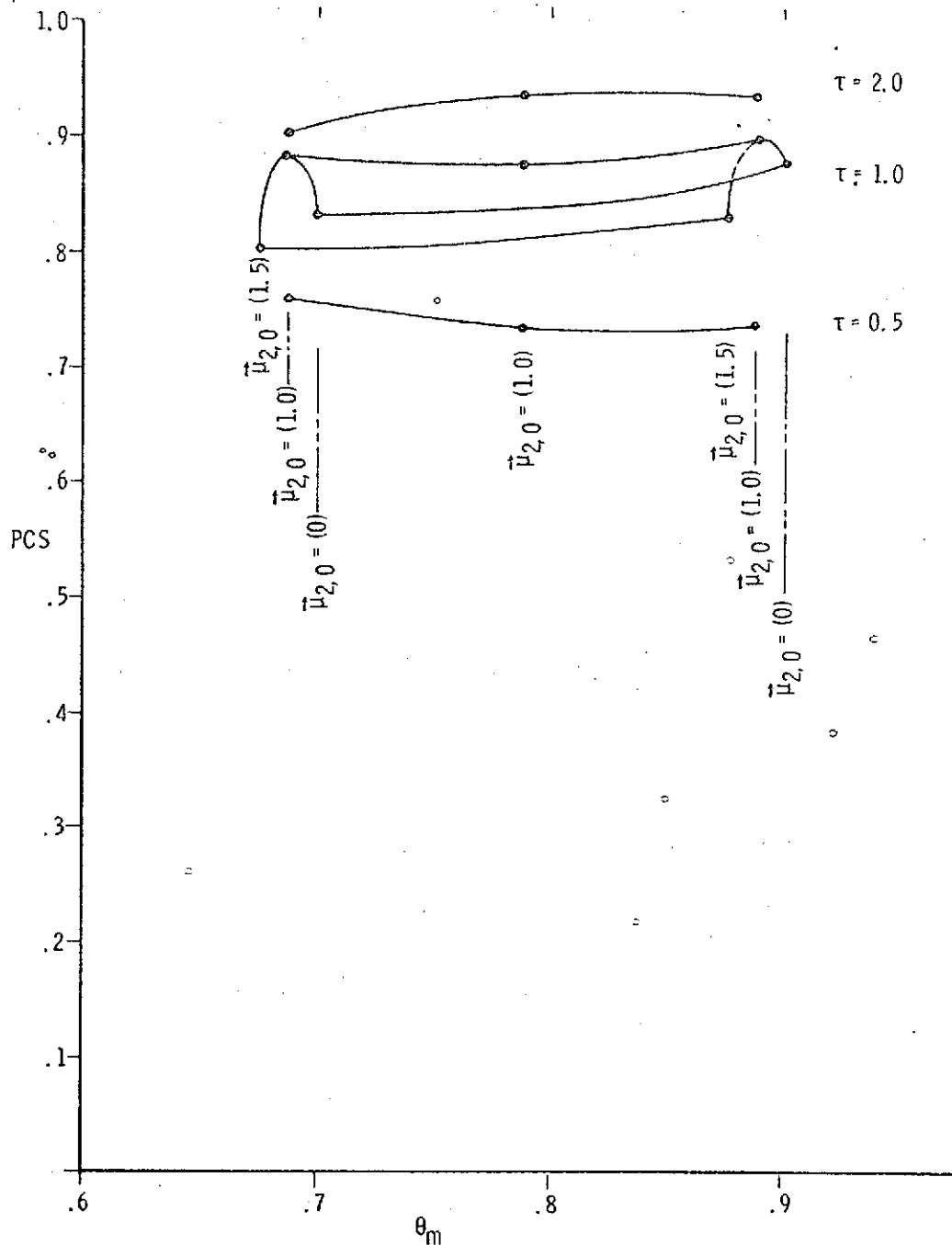


Figure 11. - Probability of correct selection (PCS) as a function of θ_m , τ , $\vec{\mu}_{2,0}$ for $J_{\max} = 8$ and H_2 true. Small sample performance simulation experiment two.

6.3.3 A Four Model Problem

In this section we study the ability of the sequential procedure to choose the correct model from the following set of completely nested model equations.

$$H_1: y = \beta_0 + \epsilon$$

$$H_2: y = \beta_0 + \beta_1 x_1 + \epsilon$$

$$H_3: y = \beta_0 + \beta_1 x_1 + \beta_2 x_2 + \epsilon$$

$$H_4: y = \beta_0 + \beta_1 x_1 + \beta_2 x_2 + \beta_3 x_1 x_2 + \epsilon$$

The prior distributions are defined by

$$\theta_{k,0} = \frac{1}{4}$$

$$\psi_{k,0} = I$$

$$\vec{\mu}_{1,0} = (0), \vec{\mu}_{2,0} = \begin{pmatrix} 0 \\ 1 \end{pmatrix}, \vec{\mu}_{4,0} = \begin{pmatrix} 0 \\ 1 \\ -1 \\ 0 \end{pmatrix}$$

And

$$\vec{\mu}_{3,0} = \begin{pmatrix} 0 \\ 1 \\ -1 \end{pmatrix}, \begin{pmatrix} 1 \\ 0.5 \\ 0 \end{pmatrix}, \begin{pmatrix} 0 \\ 0.5 \\ 0 \end{pmatrix}, \begin{pmatrix} 0 \\ 0 \\ 0.5 \end{pmatrix}$$

The equation used to generate the observations was that of H_3 with values of the parameters given by

$$\vec{\mu}^* = \begin{pmatrix} 0 \\ 1 \\ -1 \end{pmatrix}$$

The value of τ used was 1.0.

For fixed values of β_3 , the values of the prior parameter means for H_1 and H_2 and the four prior means for H_3 can be plotted in 3-space as in figure 12. For each of the four choices of $\vec{\mu}_{3,0}$, three values of θ_m ($=0.7, 0.8, 0.9$) for $J_{MAX} = 8$ were used and the resulting PCS and ASN values for the 12 combinations are given in table 15. PCS is plotted in figure 12. In all cases, 1000 simulations of the procedure were performed.

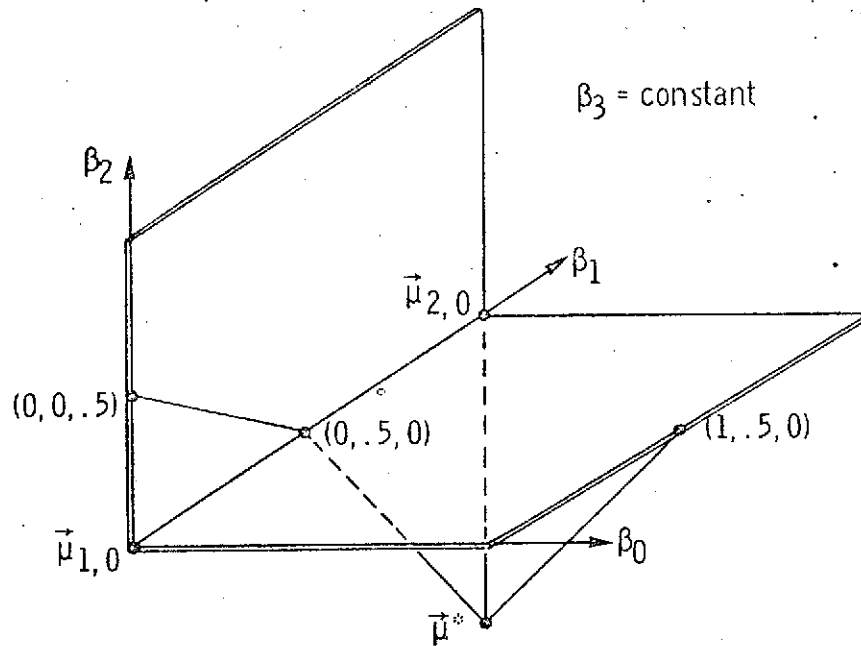


Figure 12. - Illustration of the prior means when β_3 is held constant. Small sample performance simulation experiment three.

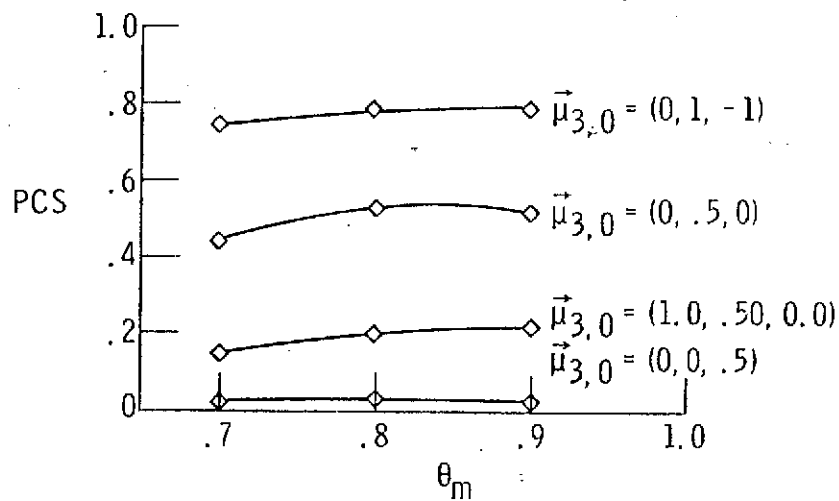


Figure 13. - Probability of correct selection (PCS) as a function of θ_m and $\vec{\mu}_{3,0}$ for $J_{\max} = 8$ and H_3 true. Small sample performance simulation experiment three.

6.4 Discussion of Results

We now make some general observations concerning the results of the simulation experiments.

First, consider the large sample results. In the context of the fact that sequential procedures are primarily developed in the hope that reliable decisions can be made with small samples rather than large samples, these results are not of primary importance. It is interesting and informative to know, however, that the procedures are consistent. Since the study of limiting posterior distributions resulting from sequentially chosen experiments is known to be an extremely difficult and delicate problem, simulation experiments may be helpful by indicating to researchers what large sample behavior is likely to be true. In the problems studied in this paper it seems quite likely that when non-nested models are encountered, the posterior probability of the true hypothesis has a limiting value of unity. It also seems most likely that the limiting posterior mean of the true hypothesis does indeed equal the values of the unknown parameters generating the data.

When nested models are encountered, however, the results are not as enlightening. It appears that if the posterior probability of the correct hypothesis does not achieve a limit of unity, it at least attains a large value and then randomly fluctuates about that value. There is indication that the conjecture of Box and Hill that for these nested models there is a distinct preference by the

sequential procedure to choose the model with the smaller number of parameters as true. For instance, the polynomial study $L = 6$, $i^* = 5$ indicates that if a model with more parameters is true but can be approximated closely by one with fewer parameters, there is a preference for the smaller model.

In examining the small sample performance simulation experiments, it is seen that PCS drops off fairly rapidly as the distance of the prior mean of the correct model from the true values of the parameters increases. This supports the conjecture of Chernoff and Meeter et al. that there may often be "initial bungling." It should be noted, however, that in all cases studied, the prior means of the competing models were all set to be as close to the true model parameter values as could be done. Thus, in a sense, these experiments can be considered to be presenting the most unfavorable situation possible to the sequential procedure. In actual application it might be more reasonable to assume that the prior distributions of all the models are mis-specified to the same extent. This problem of "initial bungling" should also indicate that the statistician should have the prior precision matrices of the parameter distributions be as vague as the prior information permits.

One approach studied by Kiefer and Sacks (1963) was to plan small initial experiments as a basis for gaining information to plan a large second experiment. An alternative not studied in this dissertation, but which seems worthy of investigation, would be to set

a lower limit, say J_{MIN} , as the minimum number of observations taken before a stopping rule is applied. The sequential procedure would use the same rule as developed for selection of experiments but large posterior probabilities on the models would be ignored until a sufficient number of observations are taken to avoid the consequences of initial bungling. This also makes sense from the point of view of obtaining parameter estimates. Surely an experimenter would not be content to terminate sampling with two or three observations even if the resulting probabilities are overwhelmingly in favor of one hypothesis unless he had extremely good prior information.

The last topic to discuss is the modified stopping and selection procedure introduced in Chapter 5. This was not applied to any of the simulation experiments performed in this dissertation. The large sample simulation results indicate that when θ_m and/or J_{MAX} are large, then this modified procedure may be of value. For the problems considered here, it is seen that even for nested models, the unmodified procedure performs quite well when J_{MAX} is small.

CHAPTER 7

EXAMPLE OF APPLICATION

This chapter first presents a general outline of the situations in which the results of this paper may be applied. Following this an example from the literature is presented. The purpose of this example is to illustrate how the information available from previous experimentation can be translated to the information required for the application of the sequential procedure developed herein.

A Bayesian framework is used in this paper because in a great many applications there does exist some prior information which can be incorporated. The Bayesian approach to statistical inference is the most natural and satisfying method of incorporating prior knowledge. This prior knowledge may arise in several ways.

For example, when expensive or large experiments are contemplated, there is often available data from pilot studies, the literature of the field, or poorly designed prior experiments. Typically, some type of regression analysis is performed on this data but there is so little data that practically no conclusions can be drawn, only recommendations for further experimentation. The resulting equations, however, provide a very convenient starting point for the application of Bayesian methods.

In other situations, an experimenter has a great deal of ex-

perience in experiments that are similar and involve factors with which he has some previous experience. In these cases it may be safe to extrapolate his acquired knowledge from the similar but different experiments to the current experiment. If so, this may indicate some characteristics of the model equation.

A third possibility might arise for example in the carrying over of laboratory results to a production process or out-of-laboratory process. In the laboratory greater control can be exerted on many variables and typically only a small number of variables may be investigated. Often one or more mechanistic models are available. When the process is taken out of the laboratory, there will be less control over other variables and they must be accounted for by adding them to the model. Thus the experimenter is faced with the situation of having a partly mechanistic model and a partly empirical model. If the mechanistic model is sufficiently smooth in the region of interest, factorial or polynomial models can be applied in these cases and prior information might indicate which interactions or terms are most likely to exist.

The example we consider is studied in Lloyd and Lipow (1962) and Draper and Smith (1966). In these books the data presented in table 16 is used to illustrate some topics in the design of experiments and multiple linear regression analysis. The dependent variable y is the chamber pressure in rocket engines put on test. The four controlled variables are

z_1 = temperature of cycle (starting)

z_2 = vibration level

z_3 = shock by dropping (temperature)

z_4 = static fire temperature

We first postulate the model equation given by equation (7-1).

$$y = \beta_0 + \beta_1 z_1 + \beta_2 z_1^2 + \beta_3 z_2 + \beta_4 z_2^2 + \beta_5 z_3 + \beta_6 z_3^2 + \beta_7 z_4 + \beta_8 z_1 z_2 \\ + \beta_9 z_1 z_3 + \beta_{10} z_1 z_4 + \beta_{11} z_2 z_3 + \beta_{12} z_2 z_4 + \beta_{13} z_3 z_4 + \varepsilon \quad (7-1)$$

The results of a multiple linear regression analysis of the model are summarized in table 17. The terms of the model are ordered in table 17 in decreasing order of descriptive significance level.

The experiment is highly saturated with respect to equation (7-1) in the sense that 14 parameters are estimated from the data from 18 distinct combinations of levels of the independent variables.

There are also quite a few high correlations among the terms of equation (7-1) and hence high correlations among the estimated parameters. The power of the resulting t-tests may be somewhat low under these circumstances. From examination of the various descriptive significance levels, the model equations tabulated in table 18a can be considered reasonable. The prior means of the distributions are also given in table 18b. The prior precision matrices may be derived by multiplying τ times the submatrices of order 3, 6, 9, and 15 of the matrix given in table 18c. How these prior distributions were determined is now described.

In multiple linear regression, under the usual normality

assumptions, the parameter estimates from the model

$$\vec{y} = M\vec{\beta} + \vec{\epsilon}$$

are given by

$$\hat{\vec{\beta}} = (M'M)^{-1}M'\vec{y}$$

and we know

$$\hat{\vec{\beta}} \sim N(\vec{\beta}, \frac{1}{\sigma^2} M'M)$$

Thus, for the first three models of table 18, the prior means and precision matrices would reasonably be the $\hat{\vec{\beta}}$ and $\frac{1}{\sigma^2} M'M$ derived by least squares analysis using the appropriate subset of data from table 16. This is how the values of $\vec{\mu}_{\ell,0}, \psi_{\ell,0}$ for $\ell = 1, 2, 3$ were derived. For $\ell = 4$ and the data of table 16, the full equation is not estimable because there are not three levels of z_4 to estimate a coefficient of z_4^2 . Thus least squares estimates were computed and $M'M$ computed for the first 14 terms of model 4. Then an essentially diffuse prior was specified with respect to β_{14} by setting the prior mean to zero and adding the last row of the matrix in table 18c to $M'M$. The diagonal term was arbitrarily chosen to make the matrix nonsingular yet not comparable to any of the other diagonal elements in magnitude.

To complete the information required, τ must be specified and $\theta_{\ell,0}$ chosen. From the data of table 17 an unbiased estimate of σ^2 is 1.85 as computed from the replicated points. Thus we may use $\tau = \frac{1}{1.85} = 0.541$. To determine the $\theta_{\ell,0}$ it will be helpful to examine the F-ratios for lack-of-fit for the first three models in

table 18. These are $F = 5.295$, $F = 2.038$, and $F = 1.871$. These statistics are significantly large at about the 0.90, 0.80, and 0.75 levels, respectively. Based upon this, the following values of $\theta_{k,0}$ seem reasonable

$$\theta_{1,0} = 0.10$$

$$\theta_{2,0} = 0.30$$

$$\theta_{3,0} = 0.30$$

$$\theta_{4,0} = 0.30$$

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APPENDIX A

COMPUTER PROGRAM FOR SIMULATION STUDIES

The general flow of operations and computations performed by the program is described by the algorithm given in Chapter 6.1. The input required to perform these computations is first described. How these computations are achieved is described briefly by giving the major functions of the subprograms constituting the complete program. A complete FORTRAN listing is given.

INPUT

The program, as presented here, can only accommodate polynomial models over the interval $[-1,+1]$ and two-level factorial models. This can be changed by writing one new subroutine (MFORM) to handle more general models. The program identifies the parameters by their integer subscripts and computes the x_{ij} values for the M-matrices according to the following convention:

1. For polynomial models, the subscript i indicates the parameter which is the coefficient of x^i .
2. For factorial models, the coefficients are assumed ordered in the standard order according to the description in Sidik and Holms (1971). The treatment combinations are also assumed to be in standard order and the independent variable values are constructed as described in Sidik and Holms (1971).
3. The order of the models as specified for input are written

such that the parameter subscripts are in increasing order.

The specific input cards are now described below and illustrated by the input for a case run in Chapter 6.3.1. The problem input is specified by the models and parameter values:

$$H_1: y = \beta_1 x_1 + \epsilon$$

$$H_2: y = \beta_2 x_2 + \epsilon$$

$$H_3: y = \beta_1 x_1 + \beta_2 x_2 + \epsilon$$

$$\left. \begin{array}{l} \psi_{\ell,0} = 1 \\ \theta_{\ell,0} = \frac{1}{3} \end{array} \right\} \ell = 1, 2, 3 \quad \begin{array}{l} \theta_m = 0.70 \\ \tau = 0.5 \end{array}$$

$$J_{\text{MAX}} = 8$$

1500 simulations

$$\vec{\mu}_{1,0} = (1.0)$$

$$\vec{\mu}_{2,0} = (1.0)$$

$$\vec{\mu}_{3,0} = \begin{pmatrix} 0.5 \\ 0.5 \end{pmatrix} \quad \vec{\mu}^* = \begin{pmatrix} 1.0 \\ 1.0 \end{pmatrix}$$

Start random number generator with 041 574 501 221.

1. IDENTIFICATION (13A6)

One card for Hollerith input description of the problem.

2. NAMELIST INPUT (\$NAML)

Most of the control parameters are included in a NAMELIST input set. The list of parameters and their purpose follows:

NFAC For factorial problems, this supplies the number of factors. For polynomial problems, this variable need not be supplied.

NHYP Number of hypothesized models.

NSYM Number of simulations to run.

MXTRTS J_{MAX}

NTREND Not used for this report. Set to zero.

LEVOUT An output control variable. Certain basic output is automatically printed. Extra intermediate output can be printed by setting LEVOUT to an integer between 1 and 7. For performance studies set to 0. For large sample runs set to 2. For debugging set to 7.

LTRUE Supplies i^* , the correct model subscript.

TAU τ

CSTOP θ_m

TMAX An upper limit on execution time. If this limit is exceeded, the program dumps for a restart.

IFSTRT Set T for supplied starting value for random number sequence. Set F if sequence is to start with initialization value. (See description of subroutines RAND and SAND for further information.)

RESTRT Is this problem a restart of a case terminated by exceeded time? T or F.

POLY T implies polynomial model. F implies factorial model.

NX For polynomial models, the x space is restricted to the interval -1 to +1. NX specifies the number of points used to discretize the interval into equal increments.

3. FORMAT FOR PRECISION MATRICES (13A6, A2)

For each set of model equations supply one set of 4A, 4B, 4C, 4D, and 4E.

4A. NUMBER OF PARAMETERS IN MODEL, PRIOR PROBABILITY OF MODEL (I6, F12.6)

4B. PARAMETER SUBSCRIPTS (13I6)

4C. PRIOR MEANS (12F6.0)

4D. TRUE VALUE OF PARAMETERS (12F6.0).

This card should be supplied only for the set corresponding to the correct model.

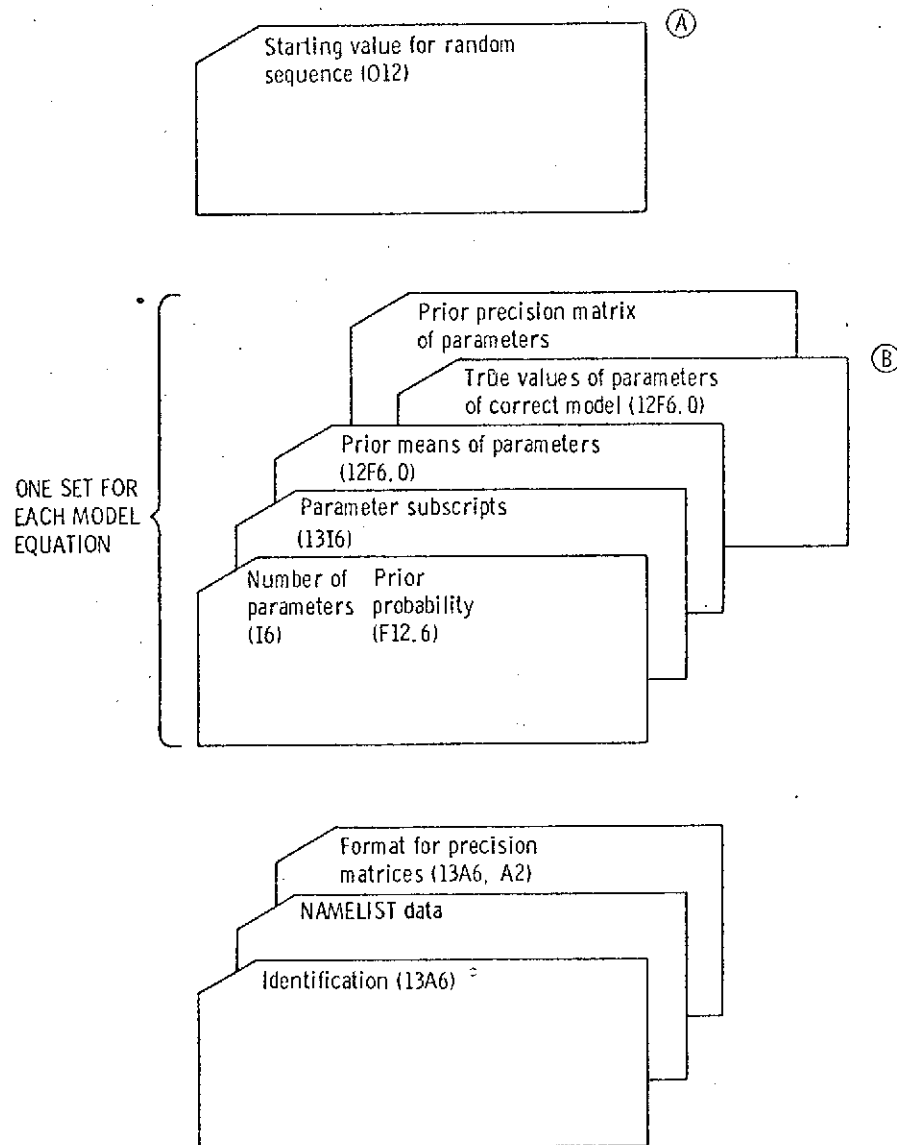
4E. PRIOR PRECISION MATRIX

Only the upper triangular half of Ψ is specified with the order being

$$\psi_{11}, \psi_{12}, \psi_{22}, \psi_{13}, \text{ etc.}$$

5. STARTING VALUE FOR PSEUDO-RANDOM SEQUENCE (012)

A graphical illustration of a data deck is given in figure A1 and a FORTRAN sheet giving the sample input is given in table A1.



A - This card is optional and its presence depends upon NAMELIST data.

B - Supplied only for true model.

Figure A-1. - Graphical illustration of card input for simulation program.

SUBPROGRAMS AND THEIR MAJOR FUNCTIONS

1. SEQDES. This is the main program. It reads the input cards, exercises general control over the other subroutines. It also outputs the final results.
2. ACOMB. Called once by SEQDES at the beginning of each case. This subroutine scans the lists of parameter subscripts for each model and constructs from them a new list of subscripts in ascending order which contains the parameters appearing in at least one model.
3. SVSTRT (GTSTRT). A double entry subprogram. The entry SVSTRT is used once at the beginning of each case to save the prior probabilities, means, and precision matrices. Then after each simulation, entry GTSTRT is used to re-initialize the working probability, mean, and precision vectors to the original values.
4. MFORM. There are two versions, MFORMX and MPOLY. The first is for factorial runs and the second for polynomial models. These subroutines construct the M-matrices required for each model when given the experiment choice (a level of x for polynomials, a treatment combination for factorials). It is called by INFX and YGEN.
5. INFX (UPDATE). Entry INFX accepts the current state of sampling as defined by $\{\theta_{\lambda,j}, \vec{u}_{\lambda,j}, \psi_{\lambda,j}; \lambda = 1, L\}$ and determines the expected K-L information for a specified experiment. Repeated calls then are made by SEQDES to determine the optimal experiment. Entry UPDATE is called by YGEN and accepts the current state values

and an observed random response to compute the posterior distributions.

6. RANDOM. Called once for each random observation taken. Generates a random observation, ε , from $N(0, \tau)$.

7. YGEN. Called once for each random observation. This routine accepts the $N(0, \tau)$ variate generated by RANDOM and calls MFORM to compute M , where M is the design matrix appropriate for the correct model and the experiment chosen as optimal. Then

$$y = M\vec{\mu}^* + \varepsilon$$

is computed. After generating y , posterior values for θ_ℓ are computed. Sampling for that simulation trial is then terminated if any θ_ℓ exceeds θ_m . If sampling is not terminated, YGEN calls UPDATE to compute the posterior $\vec{\mu}_\ell(y)$ and $\Psi_\ell(y)$, and the sampling procedure is continued.

8. COUNTX. This subroutine is called by YGEN whenever a $\theta_\ell \geq \theta_m$ or by SEQDES whenever J_{MAX} observations have been taken. It counts the number of times each model is chosen and records the distribution of sample sizes. These are then output after all simulations have been computed.

9. INVXTX. Inverts a symmetric matrix stored in the lower symmetric storage mode.

10. TRIUX. Outputs a lower triangular matrix.

11. MTVEC. Computes the product vector resulting from the multiplication of a vector and a symmetric matrix.

12. RAND (SAND). This pair of entries provide the pseudo-random sequence of uniform random variates. SAND is an initializing entry which must be called before any calls to RAND. When RAND is called, it computes the next value in the random sequence from the current value. The return argument of RAND is the floating point uniform random variate. The input and return argument of SAND is an argument which saves the integer value of the random variate. The generator is the multiplicative-congruential type obtained by taking the low order 36 bits of the product $r_{r-1} \times k$ where

r_{r-1} is the previous random number
 r_0 is 1 and
 k is 5^{15} .

The properties of this generator are discussed by Taussky and Todd (1956) and Coveyou and Macpherson (1967).

13. BCREAD (X1,X2). BCDUMP (X1,X2). These routines, respectively, read and punch cards in absolute binary at the rate of 22 words per card. The data read or punched begins at the location in core of the variable X1 and ends at the location in core of the variable X2.


```

0000100 COMMON/UNITS/ IUNIN,IUNOUT,MASK,LEVOUT,PFMT(4)
0000200 COMMON/XRAND/Y,EPS,UNIF,U
0000300 COMMON/ERR/TAU,SIGMA
0000400 COMMON/ALPH/IALPH(1000),I1ALPH(512),ALPHMU(1000),REALMU(512)
0000500 X ,NPARAM(10),MUD(11),PMU(1000)
0000600 COMMON/PS/PINDEX(11),PREC(2000)
0000700 INTEGER PINDEX
0000800 COMMON/CNTRLS/NHYP,NALL,NFAC,NFULL,NFULM1,TRIMNT,NTREND,NORS,LTRUE
0000900 INTEGER TRIMNT
0001000 COMMON/PREFNC/NAVG(1000),PCS,XMSE,XMST,INITAL ,IHCNT(10),MSOFAR
0001100 COMMON/INF/XINF
0001200 COMMON/PROBS/THETA(10),CSTOP
0001300 C
0001400 C
0001500 C
0001600 COMMON/MM/OESM(1000)
0001700 COMMON/YDIST/S(10),R(10),SV(10),RV(10)
0001800 COMMON/XTRA/AA(1000),B(2000),C(2000),D(10)
0001900 C
0002000 C** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** **
0002100 C
0002200 DATA ENDCRD/4HENDCRD/
0002300 REAL IDENT
0002400 EQUIVALENCE (XI,IX),(KS,SK)
0002500 DIMENSION IDENT(14)
0002600 DIMENSION INPFMT(14)
0002700 LOGICAL IFSTRT,RESTR
0002800 LOGICAL KPGEN
0002900 LOGICAL POLY
0003000 NAMELIST/NAML/NFAC,NHYP,NSYM,MXTRYS,NTREND,LEVOUT,LTRUE,TAU,
0003100 X CSTOP,TMAX,IFSTRT,RESTR
0003200 X ,POLY,NX
0003300 C
0003500 C*****
0003600 C
0003700 KPGEN=.FALSE.
0003800 IGOTO=1
0003900 CALL TIME1(TSTRT)
0004000 1 READ(IUNIN,5040) IDENT
0004100 IF(IDENT(1).EQ.ENDCRD) STOP
0004200 WRITE(IUNOUT,6010) IDENT
0004300 C
0004400 READ(IUNIN,NAML)
0004500 WRITE(IUNOUT,NAML)
0004600 NFULL=2**NFAC
0004700 NFULM1=NFULL-1
0004800 IF(POLY) NFULL=NX
0004900 SIGMA= 1.0/SQRT(TAU)
0005000 READ(IUNIN,5040) INPFMT
0005100 MUD(1)=1
0005200 PINDEX(1)=1
0005300 DO 190 N= 1,NHYP
0005350 IHCNT(N)=0
0005400 WRITE(IUNOUT,6365)
0005500 READ(IUNIN,5020) NPARAM(N),THETA(N)
0005700 MUD(N+1)=MUD(N)+ NPARAM(N)
0005800 PINDEX(N+1)=PINDEX(N)+ (NPARAM(N)*(NPARAM(N)+1))/2
0005900 MU0 = MUD(N)
0006000 MUHI= MUD(N+1)-1

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0006100      IPLO= PINDEX(N)
0006200      IPHI= PINDEX(N+1)-1
0006300      READ(IUNIN,5010)(IALPH(L),L=MULO,MUHI)
0006400      READ(IUNIN,5050)(ALPHMU(L),L=MULO,MUHI)
0006500      IF (N.NE.LTRUE) GO TO 100
0006600      NR = NPARAM(N)
0006700      WRITE(IUNOUT,6360)
0006800      READ(IUNIN,5050) (REALMU(L),L=1,NR)
0006900      WRITE (IUNOUT,6355) (REALMU(L),L=1,NR)
0007000      WRITE (IUNOUT,6360)
0007100 100 CONTINUE
0007200      READ(IUNIN,INPFMT)(PREC(L),L=IPLO,IPHI)
0007300      WRITE(IUNOUT,6050) N,NPARAM(N),THETA(N)
0007400      WRITE(IUNOUT,6060)(IALPH(L),ALPHMU(L),L=MULO,MUHI)
0007500      WRITE(IUNOUT,6090)
0007600      CALL TRIANG(PREC(IPLO),NPARAM(N),B,PFMT,IALPH(MULO))
0007700 C
0007800      CALL MTEVC(PREC(IPLO),ALPHMU(MULO),NPARAM(N),PMU(MULO))
0007900      IF(LEVOUT.GE.5)WRITE(IUNOUT,6060)(IALPH(L),PMU(L),L=MULO,MUHI)
0008000 190 CONTINUE
0008100 C
0008200 C
0008300 C
0008400 C
0008500 C
0008600      M = NPARAM(LTRUE)
0008700      JJ = MOD(LTRUE)-1
0008800      XMSE=0.0
0008900      DO 195 J=1,M
0009000      JJ = JJ+1
0009100 195 XMSE = XMSE + (ALPHMU(JJ)-REALMU(JJ))**2
0009200      XMST = SQRT(XMSE*TAU)
0009300      XMSE = SQRT(XMSE)
0009400      WRITE (IUNOUT,6100) XMSE,XMST
0009500      XMSE = 0.0
0009600      XMST = 0.0
0009700      DO 197 I=1,MXTRTS
0009800 197 NAVG(I) = 0
0009900      PCS = 0.0
0010000      CALL SAND(INITL)
0010100      IF(KPGEN) INITL=NTLSV
0010200      IF(1FSTRT) READ(IUNIN,5090) INITL
0010300      IF(.NOT.RESTRT) GO TO 200
0010400      WRITE(IUNOUT,6425)
0010500      CALL BCREAD(NAVG(1),MSOFAR)
0010600 C
0010700 C
0010800 C
0010900 C
0011000 C
0011100 C
0011200 C
0011300 C
0011400 C
0011500 C*****
0011600 C
0011700 200 CONTINUE
0011800      WRITE(IUNOUT,6500) INITL
0011900 6500 FORMAT(26H STARTING VALUE FOR RAND= 012)
0012000      CALL ACOMB
0012100      CALL SVSTRT
0012200 C
0012300 C*****
0012400 C
0012500      MMM=1
0012600      IF(RESTRT) MMM=MSOFAR+1
0012700      DO 800 M=MMM,NSYM

```

```

0012800      MSOFAR=M
0012900      CALL GISTRT
0013000      DO 700 NTR=1,MXTRIS
0013100      NORS= NTR
0013200      SINF= 0.0
0013300      DO 600 ITRT=1,NFULL
0013400      TRTMNT= ITRT-1
0013500      CALL MFORM
0013600      CALL INFMTH
0013700      IF (XINF.LE.SINF) GO TO 600
0013800      SINF= XINF
0013900      ISV= TRTMNT
0014000      DO 590 N=1,NHYP
0014100      SV(N)=S(N)
0014200      RV(N)=R(N)
0014300      590 CONTINUE
0014400      600 CONTINUE
0014500      TRTMNT= ISV
0014600      IF (LEVOUT.GE.3) WRITE(IUNOUT,6210) ISV,SINF
0014700      CALL RNDM
0014800      CALL YGEN(1750)
0014900      700 CONTINUE
0015000      750 CALL COUNT
0015100      CALL TIME1(TNOW)
0015200      TPRNT=(TNOW-YSTRT)/3600.
0015300      IF (TPRNT.LT.TMAX) GO TO 800
0015400      IGOFF=2
0015500      GO TO 810
0015600      800 CONTINUE
0015700 C
0015800 C
0015900 C
0016000      810 CONTINUE
0016100      WRITE(IUNOUT,6240) TPRNT
0016200      XXMSE=XMSE/FLOAT(MSOFAR)
0016300      XXMST=XXMSE*SQRT(TAU)
0016400      ASN=0.0
0016500      DO 850 I=1,MXTRIS
0016600      850 ASN=ASN+FLOAT(NAVG(I)*I)
0016700      ASN=ASN/FLOAT(MSOFAR)
0016800      WRITE(IUNOUT,6400) (NAVG(I),I=1,MXTRIS)
0016810      WRITE(IUNOUT,6600) (IHCONT(I),I=1,NHYP)
0016820      6600 FORMAT(1H 10110)
0016900      6400 FORMAT(1H 10112)
0017000      PPCS=PCSF/FLOAT(MSOFAR)
0017100      WRITE(IUNOUT,6300) ASN,PPCS,XXMSE,XXMST,INITAL
0017200      NLSV=INITAL
0017300      KPGEN=.TRUE.
0017400      GO TO (1,1000), IGOFF
0017500      1000 CALL RCDUMP(NAVG(1),MSOFAR,0)
0017600      WRITE(IUNOUT,6450) MSOFAR
0017700      STOP
0017800 C
0017900 C*****
0018000 C
0018100      5010 FORMAT(13I6)
0018200      5020 FORMAT(16,3F12.6,10I1)
0018300      5040 FORMAT(13A6,A2)
0018400      5050 FORMAT(12F6.0)

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0018500 5080 FORMAT(6I6,3F6.0,2L1)
0018600 5090 FORMAT(I012)
0018700 C
0018800 C** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** 
0018900 6010 FORMAT(11H1,13A6,1A21)
0019000 6050 FORMAT(11HKKFOR MODEL 13/16H NO. PARAMETERS 15/13H PRIOR PROB. G14.-
0019100 X5)
0019300 6060 FORMAT(15HKTHE PARAMETERS IN THE MODEL AND THEIR MEANS ARE-- // -
0019400 X (( 110,G14.5 )))
0019500 6090 FORMAT(14HKTHE PRECISION MATRIX OF THE PARAMETERS IS--1
0019600 6100 FORMAT( 49HINITIAL DISTANCE OF PRIOR MEAN FROM TRUE VALUE G14.5/-
0019700 X 18X,16HDIVIDED BY SIGMA,15X,G14.5)
0019800 6210 FORMAT(1H 5G18.8)
0019900 6240 FORMAT(24H CURRENT EXECUTION TIME F10.3)
0020000 6300 FORMAT(6HKASN= G14.2/6H PCS= G14.6/12H AVG.DIST.= G14.6/ -
0020100 X 14H NORMALIZED = G14.6/2H* 012)
0020200 6355 FORMAT(1H 8G16.8)
0020300 6360 FORMAT(1H 40(2H +))
0020400 6365 FORMAT(1H 80(1H*1))
0020500 6425 FORMAT( 37H THIS IS A RESTART OF A PREVIOUS CASE)
0020600 6450 FORMAT(55H THIS CASE WAS TERMINATED BY CLOCK. DUMPING FOR RESTART/-
0020700 X ,29H NO. SIMULATIONS COMPLETED = 16)
0020800 END

```

```

0000100 SUBROUTINE ACOMB
0000200 COMMON/UNITS/ IUNIN,IUNOUT,MASK,LEVOUT,PFMT(4)
0000300 COMMON/ALPH/IALPH(1000),I1ALPH(512),ALPHMU(1000),REALMU(512) -
0000400 X ,NPARAM(10),MUD(11),PMU(1000)
0000500 COMMON/CNTRLS/NHYP,NALL,NFAC,NFULL,NFULM1,TRTMNT,NTREND,NOBS,LTRUE
0000600 C
0000700 C*****
0000800 C
0000900 NALL = NPARAM(1)-NTREND
0001000 IF(NALL.LE.0) GO TO 110
0001100 DO 100 K =1,NALL
0001200 I1ALPH(K)= IALPH(K)
0001300 100 CONTINUE
0001400 110 IF(NHYP.LE.1) RETURN
0001500 C
0001600 DO 800 N=2,NHYP
0001700 MAXI= MUD(N+1)-1-NTREND
0001800 KK=0
0001900 K=MUD(N)-1
0002000 150 KK=KK+1
0002100 200 K=K+1
0002200 IF(KK.GT.NALL) GO TO 420
0002300 IF(K.GT.MAXI) GO TO 800
0002400 250 IF(IALPH(K)-I1ALPH(KK))300,150,400
0002500 300 KS=NALL+2
0002600 DO 350 J=KK,NALL
0002700 KS=KS-1
0002800 350 I1ALPH(KS)=I1ALPH(KS-1)
0002900 I1ALPH(KK)=IALPH(K)
0003000 NALL=NALL+1

```

```

0003100      GO TO 150
0003200      400 KK=KK+1
0003300      IF (KK=NALL) 250,250,450
0003400      420 IF (K.GT.MAXI) GO TO 400
0003500      450 NMORE=MAXI-K+1
0003600      KD=K-1
0003700      DO 500 J=1,NMORE
0003800      KS1=NALL+J
0003900      KS2=KD+J
0004000      IIALPH(KS1) = IALPH(KS2)
0004100      500 CONTINUE
0004200      NALL=NALL+NMORE
0004300      400 CONTINUE
0004400      IF (LEVOUY.GE.7)WRITE(IUNOUT,1000) NALL,(IIALPH(I),I=1,NALL)
0004500      1000 FORMAT(1H 10/(1H 25I5))
0004600      RETURN
0004700      END

```

```

0000100      SUBROUTINE SVSTRT
0000200      COMMON/ALPH/IALPH(1000),IIALPH(512),ALPHMU(1000),REALMU(512)
0000300      X ,NPARAM(10),MUD(11),PMU(1000)
0000400      COMMON/PS/PINDEX(11),PREC(2000)
0000500      INTEGER PINDEX
0000600      COMMON/CNTRLS/NHYP,NALL,NFAC,NFULL,NFULM1,TRTMNT,NJTEND,NDBS,LTRUE
0000700      COMMON/PROBS/THETA(10),CSTOP
0000800      DIMENSION XLPHMU(1000),XPREC(2000),XTHET(10),XPMU(1000)
0000900      M= MUD(NHYP+1)-1
0001000      DO 20 J=1,M
0001100      XLPHMU(J) = ALPHMU(J)
0001200      XPMU(J)=PMU(J)
0001300      20 CONTINUE
0001400      M = PINDEX(NHYP+1)-1
0001500      DO 40 J=1,M
0001600      XPREC(J) = PREC(J)
0001700      40 CONTINUE
0001800      DO 60 J=1,NHYP
0001900      XTHET(J) = THETA(J)
0002000      60 CONTINUE
0002100      RETURN
0002200      ENTRY GSTRT
0002300      M= MUD(NHYP+1)-1
0002400      DO 120 J=1,M
0002500      ALPHMU(J)=XLPHMU(J)
0002600      PMU(J)=XPMU(J)
0002700      120 CONTINUE
0002800      M=PINDEX(NHYP+1)-1
0002900      DO 140 J=1,M
0003000      PREC(J)= XPREC(J)
0003100      140 CONTINUE
0003200      DO 160 J=1,NHYP
0003300      THETA(J) = XTHET(J)
0003400      160 CONTINUE
0003500      RETURN
0003600      END

```

```

0000100 C THIS IS FOR TWO-LEVEL FACTORIALS ONLY
0000150 SUBROUTINE MFORM
0000200 COMMON/UNITS/ IUNIN,IUNOUT,MASK,LEVOUT,PFMT(4)
0000300 COMMON/ALPH/IALPH(1000),IIALPH(512),ALPHMU(1000),REALMU(512)
0000400 X ,NPARAM(10),MUD(11),PMU(1000)
0000500 INTEGER ITRMNT
0000600 COMMON/CNTRL/NHYP,NALL,NFAC,NFULL,NFULM1,TRTMNT,NTREND,NORS,LTRUE
0000700 COMMON/MM/DESM(1000)
0000800 C
0000900 C** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** **
0001000 C
0001100 DIMENSION LASTA(10)
0001200 EQUIVALENCE (KS,SK),([X,X1])
0001300 C*****
0001400 C
0001500 DO 5 N=1,NHYP
0001600 LASTA(N)=0
0001700 5 CONTINUE
0001800 C
0001900 DO 1000 I=1,NALL
0002000 IPARAM = IIALPH(I)
0002100 IF(IPARAM.NE.0) GO TO 40
0002200 A= +1.0
0002300 GO TO 500
0002400 C
0002500 40 CONTINUE
0002600 ITR=TRTMNT
0002700 DX= 1.0
0002800 DO 150 J=1,NFAC
0002900 XI=AND(MASK,IPARAM)
0003000 IX= IX+1
0003100 GO TO (130,100), IX
0003200 100 SK=AND(MASK,ITR)
0003300 KS= KS+1
0003400 GO TO (110,130), KS
0003500 110 DX= -DX
0003600 130 IPARAM= IPARAM/2
0003700 ITR=ITR/2
0003800 150 CONTINUE
0003900 A= DX
0004000 C
0004100 C
0004200 500 CONTINUE
0004300 IF(LEVOUT.GE.7) WRITE(IUNOUT,6000) IIALPH(I),A
0004400 C
0004500 DO 950 K=1,NHYP
0004600 IF(LASTA(K)-NPARAM(K) + NTREND) 520,950,950
0004700 520 IX= MUD(K)+LASTA(K)
0004800 IF(IIALPH(IX) - IALPH(IX)) 950,550,9999
0004900 550 CONTINUE
0005000 LASTA(K)= LASTA(K) +1
0005100 DESM(IX)= A
0005200 950 CONTINUE
0005300 1000 CONTINUE
0005400 IF(NTREND.LE.0) GO TO 1060
0005500 A= NOR5
0005600 R= 1.0E0
0005700 DO 1050 J=1,NTREND
0005800 B= R*A

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```

0005900      KS= J-NTREND
0006000      DO 1040 K=1,NHYP
0006100      KS=KS+NPAPAM(K)
0006200      DESM(KS)= B
0006300      1040 CONTINUE
0006400      1050 CONTINUE
0006500      1060 CONTINUE
0006600      ITR=MUD(NHYP+1)-1
0006700      IF (LEVOUT.GE.7) WRITE(IUNOUT,6010) (DESM(I),I=1,ITR)
0006800      RETURN
0006900 C
0007000      9999 CALL EXIT
0007100      RETURN
0007200      6000 FORMAT(1H 16,F6.0)
0007300      6010 FORMAT(6H DESM= / (6X,15F8.0))
0007400 C
0007500      END

```

```

0000100 C THIS IS FOR POLYNOMIALS ONLY
0000150      SUBROUTINE MFORM
0000200      COMMON/UNITS/ IUNIN,IUNOUT,MASK,LEVOUT,PFMT(4)
0000300      COMMON/ALPH/IALPH(1000),IIALPH(512),ALPHMU(1000),REALMU(512)
0000400      X ,NPAPAM(101),MUD(11),PMU(1000)
0000500      INTEGER TRTMNT
0000600      COMMON/CTRLS/NHYP,NALL,NFAC,NFULL,NFULM1,TRTMNT,NTREND,NQBS,LTRUE
0000700      COMMON/MM/DESM(1000)
0000800 C
0000900 C** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** ** **
0001000 C
0001100      DIMENSION LASTA(10)
0001200      EQUIVALENCE (KS,SK),(IX,XI)
0001300 C*****
0001400 C
0001500      DX=2.0/FLOAT(NFULL-1)
0001600      X=-1.0+FLOAT(TRTMNT)*DX
0001650      IF (ABS(X).LE.1.0E-4) X=0.0
0001700      DO 5 N=1,NHYP
0001800      LASTA(N)=0
0001900      5 CONTINUE
0002000 C
0002100      DO 1000 I=1,NALL
0002200      IPARAM = IIALPH(I)
0002300      A=X**IPARAM
0002400 C
0002500 C
0002600      500 CONTINUE
0002700      IF (LEVOUT.GE.7) WRITE(IUNOUT,6000) IIALPH(I),A
0002800 C
0002900      DO 950 K=1,NHYP
0003000      IF (LASTA(K)-NPAPAM(K) + NTREND) 520,950,950
0003100      520 IX= MUD(K)+LASTA(K)
0003200      IF (IIALPH(I) - IALPH(IX)) 950,550,9999
0003300      550 CONTINUE
0003400      LASTA(K)= LASTA(K) +1

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```

0003500      DESM(IX)= A
0003600      950 CONTINUE
0003700      1000 CONTINUE
0003800      IF(NTREND.LE.0) GO TO 1060
0003900      A= NOBS
0004000      B= 1.0E0
0004100      DO 1050 J=1,NTREND
0004200      B= B*A
0004300      KS= J-NTREND
0004400      DO 1040 K=1,NHYP
0004500      KS=KS+NPARAM(K)
0004600      DESM(KS)= B
0004700      1040 CONTINUE
0004800      1050 CONTINUE
0004900      1060 CONTINUE
0005000      ITR=MUD(NHYP+1)-1
0005100      IF (LEVOUT.GE.7) WRITE(IUNOUT,6010) (DESM(I),I=1,ITR)
0005200      RETURN
0005300 C
0005400      9999 CALL EXIT
0005500      RETURN
0005600      6000 FORMAT(1H 16,F6.0)
0005700      6010 FORMAT(6H DESM= /16X,15F8.0)
0005800 C
0005900      END

```

```

0000100      SUBROUTINE INEFTN
0000200      COMMON/UNITS/ IUNIN,IUNOUT,MASK,LEVOUT,PEFT(4)
0000300      COMMON/CNTRLS/NHYP,NALL,NFAC,NFULL,NFULM1,TRTMNT,NTREND,NOBS,LTRUE
0000400      INTEGER TRTMNT
0000500      COMMON/XRAND/Y,EPS,UNIF,U
0000600      COMMON/ERR/TAU,SIGMA
0000700      COMMON/ALPH/IALPH(1000),IIALPH(512),ALPHMU(1000),REALMU(512)
0000800      X      ,NPARAM(10),MUD(11),PMU(1000)
0000900      COMMON/PS/PINDEX(11),PREC(2000)
0001000      INTEGER PINDEX
0001100      COMMON /INF/XINF
0001200      COMMON/PROBS/THETA(10),CSTOP
0001300      COMMON/MM/DESM(1000)
0001400      COMMON/YDIST/ S(10),R(10),SV(10),RV(10)
0001500      COMMON/XTRA/AA(1000),B(2000),C(2000),D(10)
0001600      DIMENSION DESMU(10),TRACE(10)
0001700      LOGICAL UPDT
0001800 C
0001900 C*****
0002000 C
0002100      GO TO 5
0002200      ENTRY UPDATE
0002300      UPDT=.TRUE.
0002400      GO TO 6
0002500      5 UPDT=.FALSE.
0002600      6 CONTINUE
0002700      DO 10 K=1,NHYP
0002800      DESMU(K)=0.0

```



```

0002900      10 CONTINUE
0003000 C
0003100 C*****
0003200 C
0003300      DO 520 N=1,NHYP
0003400      IF(LEVOUT,GE.5) WRITE(IJNDUT,6005) N,NDBS
0003500      MS= MUD(N)
0003600      JE= NPARAM(N)
0003700 C
0003800 C
0003900      IR=0
0004000      DO 100 J=1,JE
0004100      KS1=MS+J-1
0004200      IF(UPDT) GO TO 25
0004300      DESMU(N)=DESMU(N)+DESM(KS1)*ALPHMU(KS1)
0004400      25 CONTINUE
0004500      DO 100 JJ=1,J
0004600      IR= IR+1
0004700      KS2 = MS+JJ-1
0004800      B(IR)= DESM(KS1)*DESM(KS2)
0004900      100 CONTINUE
0005000 C
0005100 C*****
0005200 C
0005300      IF(LEVOUT,GE.6) CALL TRIANG(B,JE,R,PFMT,IALPH(MS))
0005400      KS1=PINDEX(N+1)-PINDEX(N)
0005500      KS2= PINDEX(N)
0005600      DO 120 K=1,KS1
0005700      B(K)= TAU*B(K)+PREC(KS2)
0005800      C(K)=B(K)
0005900      KS2= KS2+1
0006000      120 CONTINUE
0006100      CALL INVXTX(B,JE)
0006200      IF(LEVOUT,GE.6) CALL TRIANG(R,JE,R,PFMT,IALPH(MS))
0006300 C
0006400 C** ** ** ** **
0006500      IF(.NOT.UPDT) GO TO 199
0006600      KS1=PINDEX(N+1)
0006700      IPLO = PINDEX(N)
0006800      KS1=KS1-IPLO
0006900      MUH1=MUD(N+1)-1
0007000      KS2=IPLO
0007100      DO 130 J=1,KS1
0007200      PREC(KS2)=C(J)
0007300      KS2=KS2+1
0007400      130 CONTINUE
0007500      KS1=MS-1
0007600      DO 150 J=1,JE
0007700      KS1= KS1+1
0007800      AA(J)= Y*DESM(KS1)*TAU + PMU(KS1)
0007900      150 CONTINUE
0008000      CALL MIVC(B,AA,JE,ALPHMU(MS))
0008100      CALL MIVC(PREC(IPLO),ALPHMU(MS),JE,PMU(MS))
0008200      IF(LEVOUT,LT.3) GO TO 520
0008300      IF(LEVOUT,GE.5) GO TO 160
0008400      WRITE(IJNDUT,6020) (ALPHMU(L),L=MS,MUH1)
0008500      IF(LEVOUT,LT.4) GO TO 520
0008600      WRITE(IJNDUT,6090)
0008700      CALL TRIANG(PREC(IPLO),JE,R,PFMT,IALPH(MS))

```

```

0008900      GO TO 520
0008900 160 CONTINUE
0009000      WRITE(IUNOUT,6060)((ALPH(L),ALPHMU(L),PMU(L),L=MS,MUHT)
0009100      WRITE(IUNOUT,6090)
0009200      CALL TRIANG(PREC(IPLO),JF,8,PFMT,IALPH(MS))
0009300      GO TO 520
0009400 C
0009500 C *****
0009600 C
0009700 C
0009800 199 CALL MTVEC(B,DESM(MS),JE,AA)
0009900      DO 300 KS2=1,JE
0010000 280 AA(KS2)= TAU*AA(KS2)
0010100 300 CONTINUE
0010200      IF(LEVOUT.GE.5) WRITE(IUNOUT,60101(AA(J),J=1,JE)
0010300 C
0010400 C** ** ** ** **
0010500 C
0010600      R(N)= 0.0
0010700      MMS=MS-1
0010800      DO 360 J=1,JE
0010900      MMS= MMS+1
0011000      R(N)= R(N) + AA(J)*DESM(MMS)
0011100 360 CONTINUE
0011200      R(N)= TAU*(1.0-R(N))
0011300 C
0011400 C** ** ** ** **
0011500 C
0011600      MMS=MS
0011700      C=0.0
0011800      DO 500 J=1,JE
0011900      C=C+AA(J)*PMU(MMS)
0012000      MMS= MMS+1
0012100 500 CONTINUE
0012200      S(N)=C/R(N)
0012300      IF(LEVOUT.GE.5) WRITE(IUNOUT,6020) R(N),S(N)
0012400 520 CONTINUE
0012500 C
0012600      IF(UPDT) RETURN
0012700 C** ** ** **~
0012800 C
0012900      DO 1000 N=1,NHYP
0013000      C= 0.0
0013100      DO 580 J=1,NHYP
0013200      IF(N.EQ.J) GO TO 580
0013300      C=C+THETA(J) *P(J)
0013400 580 CONTINUE
0013500      TRACE(N)=C/R(N)
0013600      IF(LEVOUT.GE.5) WRITE(IUNOUT,6035) TRACE(N)
0013700 1000 CONTINUE
0013800 C
0013900 C*****
0014000 C
0014100      XINF = THETA(1) * TRACE(1)
0014200      DO 1500 N=2,NHYP
0014300      JE= N-1
0014400      XINF = XINF + THETA(N) * TRACE(N)
0014500      DO 1450 J=1,JE
0014600      C= S(N) - S(J)

```

```

0014700      R=C*(R(NI)+R(J))
0014800      XINF = XINF + THETA(N) *THE1A(N) *R
0014900      1450 CONTINUE
0015000      1500 CONTINUE
0015100      IF (LEVOUT.GE.5) WRITE(IUNOUT,6020) XINF
0015200      RETURN
0015300 C
0015400 C      *****
0015500 C
0015600 C
0015700 C*****
0015800 C      FORMATS
0015900 C
0016000      6005 FORMAT(3H N= I6,7H NOBS= I6)
0016100      6010 FORMAT(2H A/(1H 10G12.4))
0016200      6020 FORMAT(1H 10G12.4)
0016300      6035 FORMAT(7H TRACE= G12.4)
0016400      6060 FORMAT(51H THE PARAMETERS IN THE MODEL, AND THEIR MEANS ARE-- //
0016500      X ((1H I10,2G14.5))
0016600      6090 FORMAT(44H THE PRECISION MATRIX OF THE PARAMETERS IS-- )
0016700      END

```

```

0000100      SUBROUTINE RNDM
0000200      COMMON/XRAND/ Y,EPS,UNIF,U
0000300      COMMON/UNITS/ IUNIN,IUNOUT,MARK,LEVOUT
0000400      COMMON/ERR/TAU,SIGMA
0000500      DATA A0/2.515517/,A1/.8028577/,A2/.0103487/,
0000600      X R1/1.432788/,R2/.1892697/,R3/.0013087/,
0000700      X UMIN/.0000001/
0000800      U=U
0000900      CALL RANDI(UNIF)
0001000      U=UNIF
0001100      IF(UNIF.GT..50) U=1.0-UNIF
0001200      IF(U.LT.UMIN) GO TO 100
0001300      T2=ALOG(1.0/(U*U))
0001400      T=SQRT(T2)
0001500      EPS=T-(A0+A1*T+A2*T2)/(1.0+R1+T+R2*T2+R3*T2)
0001600      IF (LEVOUT.GE.7) WRITE(IUNOUT,6000) U,T,T2,EPS
0001700      IF(UNIF.LT..50) EPS=-EPS
0001800      EPS=EPS*SIGMA
0001900      RETURN
0002000      100 EPS=-1.0E15
0002100      IF(UNIF.GT..50) EPS=-EPS
0002200      RETURN
0002300      6000 FORMAT(6H RAND /1H 5G16.8)
0002400      END

```

```

0000100      SUBROUTINE YGEN(*)
0000200      COMMON/UNITS/IUNIN,IUNOUT,MASK,LEVOUT,PFMT(4)
0000300      COMMON/ALPH/IALPH(1000),IFALPH(512),ALPHMU(1000),REALMU(512)
0000400      X      ,NPARAM(10),MUD(111),PMU(1000)
0000500      COMMON/CNTRL/S/NHYP,NALL,NFAC,NFULL,NFULM1,TRTMNT,NTREND,NOBS,LTRUE
0000600      INTEGER TRTMNT
0000700      COMMON/MM/DESM(1000)
0000800      COMMON/XRAND/Y,EPS,UNIF,U
0000900      COMMON/YOIST/S(10),R(10),SV(10),RV(10)
0001000      COMMON/PROBS/THETA(10),CSTOP
0001100      COMMON/PRFMNC/NAVG(1000),PCS,XMSE,XMST,INITAL ,IHCNT(10),MSQFAR
0001400      EQUIVALENCE (SK,KS1,(IX,XI))
0001600 C
0001700 C
0001800      CALL MFORM
0001900      Y=EPS
0002000      M = NPARAM(LTRUE)
0002100      JJ = MUD(LTRUE)-1
0002200      DO 550 J=1,M
0002300      JJ = JJ +1
0002400      Y = Y + DESM(JJ)*REALMU(J)
0002500 550 CONTINUE
0002600      SUM = 0.0
0002700      DO 700 N=1,NHYP
0002800      C=Y-SV(N)
0002900      Q=RV(N)*C*C
0003000      S(N)=SV(N)
0003100      R(N)=RV(N)
0003200      Q=-.50*Q
0003300      IF(ABS(Q)-60.0) 660,670,670
0003400 660 IF(THETA(N)-1.0E-18) 670,670,680
0003500 670 THETA(N) = 0.0
0003600      GO TO 700
0003700 680 A = ALOG(THETA(N)) + Q
0003800      IF(ABS(A).GE.70.0) GOTO 670
0003900 690 THETA(N) = THETA(N) * SQRT(R(N)) * EXP(Q)
0004000      SUM = SUM + THETA(N)
0004100 700 CONTINUE
0004200      I=0
0004300      DO 710 N=1,NHYP
0004400 710 THETA(N)=THETA(N)/SUM
0004500      DO 920 II=1,NHYP
0004600      IF(THETA(II)-CSTOP) 920,910,910
0004700 910 I=1
0004800      GO TO 1000
0004900 920 CONTINUE
0005000 1000 CONTINUE
0005100      IF(LEVOUT.GE.2) WRITE(IUNOUT,6010) TRTMNT,(THETA(N),N=1,NHYP)
0005200 6010 FORMAT(1H 110,10F10.4)
0005300      IF(LEVOUT.GE.5) WRITE(IUNOUT,6000) Y,EPS,(S(N),R(N),THETA(N),C(N),-
0005400      X      N = 1, NHYP)
0005500 6000 FORMAT( 3H Y=G12.5,5H EPS=G12.5, 12H S,R,THETA,C /
0005600      X (2X,4G16.5))
0005700      CALL UPDATE
0005800      IF(I.EQ.1) RETURN 1
0005900      RETURN
0006000      END

```

```

0000100      SUBROUTINE COUNT
0000200      COMMON/UNITS/ IUNIN, IUNOUT, MASK, LEVOUT
0000300      COMMON/ALPH/IALPH(1000), IIALPH(512), ALPHMU(1000), REALMU(512)
0000400      X      , NPARAM(10), MUD(11), PMU(1000)
0000500      COMMON/CNTRL/NHYP, NALL, NFAC, NFULL, NFULM, TRTMT, NTREND, NOBS, LTRUE
0000600      COMMON/PROBS/THETA(10), CSTOP
0001000      COMMON/PRFMNC/NAVG(1000), PCS, XMSE, XMST, INITAL, IHCNT(10), MSOFAR
0001100      COMMON/ERR/TAU, SIGMA
0001300      XM = 0.0
0001400      NAVG(NOBS) = NAVG(NOBS) + 1
0001500      M = NPARAM(LTRUE)
0001600      JJ = MUD(LTRUE)-1
0001700      DO 5 J=1, M
0001800      JJ = JJ + 1
0001900      XM = XM + (ALPHMU(JJ) - REALMU(JJ))**2
0002000      5 CONTINUE
0002100      XMSE = SQRT(XM) + XMSE
0002200      XMST = SQRT(XM*TAU)
0002300      XM = SQRT(XM)
0002400      IF (LEVOUT.GT.1) WRITE(IUNOUT,6000) NOBS, XM, XMST
0002500      IMIN = 0
0004200      CMAX=0.0
0004300      DO 910 I=1, NHYP
0004400      IF (THETA(I) - CMAX) 910, 910, 905
0004500      905 CMAX=THETA(I)
0004600      IMIN=I
0004700      910 CONTINUE
0004800      945 CONTINUE
0004900      IF (IMIN.EQ.LTRUE) PCS = PCS + 1.0
0004950      IHCNT(IMIN)=IHCNT(IMIN)+1
0005000      IF (LEVOUT.LT.2) RETURN
0005100      DO 950 N=1, NHYP
0005200      MU0 = MUD(N)
0005300      MUH1 = MUD(N+1)-1
0005400      WRITE (IUNOUT,6010) THETA(N), (ALPHMU(L), L=MU0, MUH1)
0005500      950 CONTINUE
0005600      RETURN
0005700      6000 FORMAT(22H *****OBSERVATION NO. 14,7H XMSE= G11.3,7H XMST= G11.3 -
0005800      X ,6H***** )
0005900      6010 FORMAT(1H F10.6 , 10G12.4/(27X10G12.4))
0006000      END

```

```

0000100      SUBROUTINE INVXTX(A,NN)
0000200 C      ASSUMES THE MATRIX A IS SYMMETRIC AND POSITIVE DEFINITE, AND ONLY
0000300 C      THE UPPER TRIANGLE IS STORED AS A ONE-DIMENSIONAL ARRAY IN THE
0000400 C      ORDER A(1,1), A(1,2), A(2,2), A(1,3), A(2,3), A(3,3), ..., A(N,N).
0000500 C      NN IS THE ORDER N OF THE INPUT MATRIX A.
0000700 C
0000800      DIMENSION A(1)
0001500      D=1.0
0001600      N = NN
0001700      ITRI = 0
0001900      DO 145 K=1, N
0002000 C
0002100      ITRI = ITRI+K-1

```

```

0002200      KP1 = K+1
0002300      KM1 = K-1
0002400      KK = ITR1+K
0002600      PV = 1.000/A(KK)
0002700 C
0002800      ITR2 = 0
0002900      IF (K-1) 150,80,50
0003000 C
0003100 C      REDUCE TOP PART OF TRIANGLE, LEFT OF PIVOTAL COLUMN
0003200 50 DO 60 J=1,KM1
0003300      ITR2 = ITR2+J-1
0003400      KJ = ITR1+J
0003500      F = A(KJ)*PV
0003600      DO 60 I=1,J
0003700      IJ = ITR2+I
0003800      IK = ITR1 + I
0003900 60 A(IJ) = A(IJ) + A(IK)*F
0004000 C
0004100      IF (K-N) 70,120,150
0004200 C
0004300 C      REDUCE REST OF TRIANGLE, RIGHT OF PIVOTAL COLUMN
0004400 70 ITR2 = ITR1
0004500 80 DO 110 J=KP1,N
0004600      ITR3 = ITR1
0004700      ITR2 = ITR2+J-1
0004800      KJ = ITR2+K
0004900      F = A(KJ)*PV
0005000      DO 100 I=1,J
0005100      IF (I-K) 90,100,95
0005200 90 IJ = ITR2+I
0005300      IK = ITR1 + I
0005400      A(IJ) = A(IJ) - A(IK)*F
0005500      GO TO 100
0005600 95 IJ = ITR2 + I
0005700      ITR3 = ITR3 + I - 1
0005800      IK = ITR3 + K
0005900      A(IJ) = A(IJ) - A(IK)*F
0006000 100 CONTINUE
0006100 110 CONTINUE
0006200 C
0006300 C      DIVIDE PIVOTAL ROW-COLUMN BY PIVOT, INCLUDING APPROPRIATE SIGNS
0006400 120 ITR2 = ITR1
0006500      DO 140 I=1,N
0006600      IF (I-K) 125,130,135
0006700 125 IK = ITR1+I
0006800      A(IK) = -A(IK)*PV
0006900      GO TO 140
0007000 C      (REPLACE PIVOT BY RECIPROCAL)
0007100 130 A(KK) = PV
0007200      GO TO 140
0007300 135 ITR2 = ITR2+I-1
0007400      KI = ITR2+K
0007500      A(KI) = A(KI)*PV
0007600 140 CONTINUE
0007700 C
0007800 145 CONTINUE
0007900 C
0008000 150 RETURN
0008100      END

```

```

0000100      SUBROUTINE TRIANG(A,NN,NCOL,FORMAT,IDOUT)
0000200      DIMENSION A(1),FORMAT(1),IDOUT(1)
0000300      COMMON/UNITS/IUNIN,IUNOUT
0000400      1 FORMAT(1HK)
0000500      N = NN
0000600      NCOL = NCOL
0000700      KLUMPS = N/NCOL
0000800 C
0000900      KEEPTR = 0
0001000      K1 = 1
0001100      K2 = NCOL - 1
0001200      K3 = NCOL
0001300      IF (KLUMPS .EQ. 0) GO TO 120
0001400 C
0001500      DO 90 KLUMP=1,KLUMPS
0001600      ITR1 = KEEPTR
0001700      I = -1
0001800      ILO = (KLUMP-1)*NCOL + ITR1 + 1
0001900      DO 30 K=K1,K2
0002000      I = I + 1
0002100      ITR1 = ITR1 + K - 1
0002200      ILO = ILO + K - 1
0002300      IHI = ILO + 1
0002400      30 WRITE(IUNOUT,FORMAT) IDOUT(K),(A(J),J=ILO,IHI)
0002500      KEEPTR = ITR1 + K2
0002600      DO 60 K=K3,N
0002700      ITR1 = ITR1 + K - 1
0002800      ILO = ILO + K - 1
0002900      IHI = ILO + NCOL - 1
0003000      60 WRITE(IUNOUT,FORMAT) IDOUT(K),(A(J),J=ILO,IHI)
0003100      K1 = K1 + NCOL
0003200      K2 = K2 + NCOL
0003300      K3 = K3 + NCOL
0003400      90 WRITE(IUNOUT,1)
0003500 C
0003600      120 ITR1 = KEEPTR
0003700      IF (K) .GT. N) GO TO 180
0003800      I = -1
0003900      ILO = KLUMPS*NCOL + ITR1 + 1
0004000      DO 150 K=K1,N
0004100      I = I + 1
0004200      ITR1 = ITR1 + K - 1
0004300      ILO = ILO + K - 1
0004400      IHI = ILO + 1
0004500      150 WRITE(IUNOUT,FORMAT) IDOUT(K),(A(J),J=ILO,IHI)
0004600 C
0004700      180 RETURN
0004800      END

```

```
0000100      SUBROUTINE MTVEC(A,R,NN,C)
0000200      DIMENSION A(1),R(1),C(1)
0000300      N=NN
0000400      DO 500 J=1,N
0000500      C(J)=0.0
0000600      KADD=((J-1)*J)/2
0000700      DO 200 K=1,J
0000800      KADD=KADD+1
0000900      C(J)=C(J)+A(KADD)*R(K)
0001000 200 CONTINUE
0001100      K1= J+1
0001200      IF(K1-N) 250,250,500
0001300 250 DO 300 K=K1,N
0001400      KADD= KADD+K-1
0001500      C(J)= C(J)+ A(KADD)*R(K)
0001600 300 CONTINUE
0001700 500 CONTINUE
0001800      RETURN
0001900      END
```


APPENDIX B

LIST OF SYMBOLS

A	the space of allowable experiments
A_j	the space of allowable experiments requiring exactly j observations
a	element of A
$a^{(i)}$	the i^{th} experiment in A
a_j	experiment in A performed at the j^{th} stage of sampling
\vec{B}	vector of parameters appearing in combined model equations
$E\{X\}$	expectation of the random variable X
$E\{X Y\}$	conditional expectation of the random variable X given the value of Y
$\mathcal{E}(w)$	entropy of the probabilities at state w
$\mathcal{E}[w(\vec{y}), a]$	entropy of the posterior probabilities if system is in state w and the value \vec{y} is observed
F	element of \mathfrak{F}
\mathfrak{F}	sigma field of Borel sets over Ω
$f_\ell(\vec{y} a, \vec{\alpha})$	density function of \vec{y} under model ℓ when $\vec{\alpha}$ is given and experiment a is to be performed.
$f_\ell(\vec{y} a)$	marginal density function of \vec{y} under model ℓ when experiment a is to be performed.

$f_{\ell}(\vec{y}_{j+1} a, \vec{\alpha})$	density function of \vec{y}_{j+1} under model ℓ when a and $\vec{\alpha}$ are given
$f_{\ell}(\vec{y}_{j+1} a)$	marginal density function of \vec{y}_{j+1} under model ℓ when a is given
$g_i(w w_{i-1})$	density function of w given w_{i-1}
H_{ℓ}	denotes hypothesis ℓ about the form of the model equation
$h_i(z_1, \dots, z_k)$	function of controlled variables defining x_i
${}_0I(w, a)$	expected information in experiment a when state of system is w
$I(w, a, i, j)$	expected information for discriminating in favor of H_i against H_j in experiment a when state of system is w
\hat{i}	denotes subscript of hypothetical model with largest posterior probability
i^*	true model equation number
J_{MAX}	upper limit on total number of observations
$k(j)$	superscript of experiment performed at stage j of sampling
L	number of model equations or hypotheses postulated
M	design matrix
M_{ℓ}	design matrix for model ℓ
$N(A)$	number of elements in A
$N(\vec{\mu}, T)$	normal distribution with mean vector $\vec{\mu}$ and precision matrix T

\vec{n}	vector of n_i
n_i	number of observations taken at stage i
$n(i,j)$	number of times experiment $a^{(i)}$ is performed in j stages
$\Pr\{X\}$	probability of event X
$\Pr\{X Y\}$	probability of event X given event Y
P_i	limiting proportion of times $a^{(i)}$ performed in an infinite sequence of experiments
Q, Q_1, Q_2	denote quadratic forms
$R_{\ell,j}$	precision matrix of distribution of y_j under model ℓ
$R(w,a)$	expected reduction in entropy if experiment a is performed and state is w
$\vec{s}_{\ell,j}$	mean vector of distribution of y_j under model ℓ
T	precision matrix of distribution of ε
W	random variable defined over Ω
w	element of Ω . An observed value
X	vector of x_i
x_i	value of $h_i(z_1, \dots, z_k)$
$x_{i,j}$	value of $h_i(z_1, \dots, z_k)$ at j^{th}
\vec{y}, y	observed variable
z_i	controlled variable i
$\vec{\alpha}_{\ell}$	vector of parameters in model equation ℓ
β_i	coefficient of x_i
$\beta_i^{(\ell)}$	coefficient of x_i in model ℓ

Y	defined on page 41
δ_{ij}	Kronecker delta function
$\vec{\varepsilon}$	vector of observation errors
θ_i	probability model i is correct
$\theta_{i,j}$	posterior probability that model i is correct after j stages of sampling
θ_m	stopping probability
$\vec{\mu}_\ell$	mean vector of distribution of parameters in model ℓ .
$\xi_{\ell,j}(\vec{\alpha})$	density function of parameters in model ℓ after j stages of sampling
τ	precision of distribution of ε
$\Psi_{\ell,j}$	precision matrix of distribution of parameters of model ℓ after j stages of sampling
Ω	state space of process. Defined on page 19
$\Omega \times Y$	direct product of state space and observed variable space
$\vec{0}$	vector of zeros
\propto	proportional to
$ $	determinant of a matrix
\approx	approximately equal to
\sim	distributed as
\supset	includes

APPENDIX C

TABLES

TABLE 1. - SUMMARY OF SIMULATION RESULTS PRESENTED IN TABLES 2 THROUGH 9

Model number	Parameter	L = 4, i* = 2		L = 4, i* = 3		L = 6, i* = 3		L = 6, i* = 5	
		100 obs	500 obs	100 obs	500 obs	100 obs	500 obs	100 obs	500 obs
1	θ_1	0	0	0	0	0	0	0	0
2	θ_2	.966	.983	0	0	0	0	.010	0
3	θ_3	.032	.016	.922	.962	.860	.877	.902	.941
4	θ_4	.002	.001	.078	.038	.109	.107	.062	.023
5	θ_5					.022	.012	.017	.029
6	θ_6					.009	.004	.009	.006
1	β_0	0.0971	0.0891	0.1322	0.1377	0.1116	0.1337	0.0278	-0.0240
2	β_0	0.1010	0.0951	0.1313	0.1384	0.1253	0.1286	0.0316	0.0382
	β_1	.4981	.5025	.2485	.2522	.2411	.2544	.5114	.4977
3	β_0	0.1013	0.0941	-0.0070	-0.0067	-0.0038	-0.0015	-0.0249	-0.0099
	β_1	.4981	.5025	.2478	.2525	.2493	.2505	.5086	.5032
	β_2	-.0007	.0021	.2578	.2634	.2157	.2544	.1179	.0981
4	β_0	0.1010	0.0943	-0.0070	-0.0335	-0.0035	-0.0015	-0.0237	-0.0097
	β_1	.5029	.4915	.2497	.2416	.2572	.2429	.5168	.5046
	β_2	-.0004	.0018	.2579	.2634	.2645	.2544	.1168	.0976
	β_3	-.0076	.0112	-.0026	.0146	-.0089	.0102	-.0083	-.0017
5	β_0					0.0040	0.0015	-0.0157	0.0001
	β_1					.2564	.2440	.5120	.5009
	β_2					.2222	.2577	-.0020	.0227
	β_3					.0032	.0088	-.0048	.0024
	β_4					.0360	-.0015	.0908	.0730
6	β_0					0.0030	-0.0037	-0.0180	-0.0021
	β_1					.2390	.2559	.5057	.4638
	β_2					.2291	.2619	.0603	.0376
	β_3					.0596	-.0373	.0438	.1340
	β_4					.0301	-.0050	.0553	.0598
	β_5					-.0534	.0345	-.0413	-.0969

The column headings give the values of L and i* and the number of observations. The row headings present the parameters whose average posterior values are given. The probabilities listed for 100 observations are the averages after five simulations of 100 observations and the values after the first 100 observations of the 500 observation simulations. The averages of the posterior parameter means are based only upon the five full simulations of 100 and 500 observations, respectively. The posterior probabilities for 500 observations are based upon five simulations of 500 observations each.

TABLE 2. - $L = 4$, $i^* = 2$

Model	Param	After 100 observations					After first 100 of 500 observations				
1	θ_1	0	0	0	0	0	0	0	0	0	0
2	θ_2	.973	.979	.974	.976	.975	.976	.976	.931	.977	.923
3	θ_3	.025	.019	.024	.023	.024	.023	.023	.063	.022	.071
4	θ_4	.002	.001	.002	.001	.002	.001	.001	.006	.001	.006
1	β_0	0.0795	0.1017	0.0906	0.1271	0.0865	*	*	*	*	*
2	β_0	0.1187	0.1017	0.0753	0.1032	0.1059	*	*	*	*	*
	β_1	.5192	.5022	.4935	.4892	.4865					
3	β_0	0.1263	0.1021	0.0682	0.1067	0.1033	*	*	*	*	*
	β_1	.5191	.5021	.4935	.4894	.4866					
	β_2	-.0152	-.0008	.0141	-.0069	.0055					
4	β_0	0.1239	0.1022	0.0688	0.1059	0.1041	*	*	*	*	*
	β_1	.4858	.5044	.4771	.5186	.5288					
	β_2	-.0126	-.0010	.0133	-.0051	.0032					
	β_3	.0351	-.0024	.0170	-.0350	-.0527					
	p_0	0.25	0.23	0.23	0.17	0.19	0.18	0.23	0.27	0.21	0.25
	p_1	0	0	0	0	0	0	0	0	0	0
	p_2	.05	.05	.03	.17	.20	.24	.06	.04	.02	.01
	p_3	.02	.03	.02	.07	0	.02	.01	.01	.03	0
	p_4	.35	.20	.14	.01	.17	.04	.26	.43	.11	.25
	p_5	.06	.20	.27	.14	.06	.12	.14	.01	.19	.24
	p_6	.01	0	.06	.09	.02	.03	.01	.01	.17	0
	p_7	.05	.05	0	.13	.21	.19	.06	.02	.01	.01
	p_8	0	0	0	0	0	0	0	0	0	0
	p_9	.21	.24	.25	.22	.15	.19	.23	.21	.26	.24

* Not recorded.

The values of the posterior probabilities and parameter means after ten simulations, of 100 observations each, of the sequential selection procedure. The last five columns are data from the first 100 observations of the 500 observation simulations tabulated in table 3. The posterior means were not recorded for these cases. Also listed are the proportions p_i of the times each $a^{(i)}$ was chosen as the optimal experiment.

TABLE 3. - $L = 4$, $i^* = 2$

Model	Param	After 500 observations				
1	θ_1	0	0	0	0	0
2	θ_2	.991	.985	.990	.991	.957
3	θ_3	.009	.015	.009	.009	.040
4	θ_4	0	0	0	0	.003
1	β_0	0.0875	0.0599	0.0905	0.0757	0.1317
2	β_0	0.0921	0.0984	0.0999	0.0942	0.0909
	β_1	.4964	.5010	.5028	.5014	.5108
3	β_0	0.0923	0.1032	0.0984	0.0937	0.0827
	β_1	.4964	.5010	.5028	.5014	.5108
	β_2	-.0005	-.0096	.0029	.0012	.0163
4	β_0	0.0923	0.1022	0.0985	0.0935	0.0850
	β_1	.4948	.4886	.5043	.4882	.4817
	β_2	-.0004	-.0086	.0028	.0014	.0139
	β_3	.0016	.0126	-.0017	.0141	.0293
	P_0	0.234	0.264	0.236	0.236	0.228
	P_1	0	0	0	0	.002
	P_2	.050	.010	.056	.044	0
	P_3	.008	0	.004	.008	0
	P_4	.280	.422	.302	.306	.076
	P_5	.110	.060	.070	.088	.418
	P_6	.018	0	.026	.076	.006
	P_7	.072	.010	.088	.038	.002
	P_8	0	0	0	0	0
	P_9	.228	.226	.218	.204	.268

The values of the posterior probabilities and parameter means after 5 simulations, of 500 observations each, of the sequential selection procedure. Also listed are the proportions p_i of the times each $a^{(i)}$ was chosen as the optimal experiment.

TABLE 4. - $L = 4, i^* = 3$

Model	Param	After 100 observations					After first 100 of 500 observations				
1	θ_1	0	0	0	0	0	0	0	0	0	0
2	θ_2	0	0	0	0	0	0	0	0	0	0
3	θ_3	.788	.828	.967	.966	.966	.962	.916	.966	.941	.920
4	θ_4	.212	.172	.033	.034	.034	.038	.084	.034	.059	.080
1	θ_0	0.1140	0.1548	0.1260	0.1292	0.1368	*	*	*	*	*
2	θ_0	0.1183	0.1515	0.1255	0.1286	0.1324	*	*	*	*	*
	θ_1	.2452	.2533	.2385	.2575	.2482					
3	θ_0	-0.0043	-0.0089	-0.0159	0.0104	-0.0162	*	*	*	*	*
	θ_1	.2467	.2510	.2389	.2552	.2470					
	θ_2	.2263	.2939	.2683	.2216	.2788					
4	θ_0	-0.0045	-0.0091	-0.0158	0.0105	-0.0162	*	*	*	*	*
	θ_1	.1838	.3095	.2395	.2633	.2523					
	θ_2	.2268	.2945	.2681	.2215	.2788					
	θ_3	.0833	-.0779	-.0009	-.0107	-.0070					
	P_0	0.18	0.17	0.17	0.17	0.17	0.18	0.17	0.17	0.17	0.17
	P_1	0	0	0	0	0	0	0	0	0	0
	P_2	.32	.31	.30	.30	.29	.32	.30	.30	.31	.30
	P_3	0	.01	0	.02	0	0	0	0	.01	.02
	P_4	0	0	.03	0	.02	0	.03	.02	0	0
	P_5	.03	.01	.02	.03	.04	.02	.02	.02	.01	.05
	P_6	0	0	.01	.02	0	.01	0	0	0	.02
	P_7	.30	.32	.30	.28	.30	.30	.31	.31	.32	.28
	P_8	0	0	0	0	0	0	0	0	0	0
	P_9	.17	.18	.17	.18	.18	.17	.17	.18	.18	.16

The values of the posterior probabilities and parameter means after 10 simulations, of 100 observations each, of the sequential selection procedure. The last 5 columns are data from the first 100 observations of the 500 observation simulations tabulated in table 5. The posterior means for these 5 cases were not tabulated. Also listed are the proportions p_i of the times each $a(i)$ was chosen as the optimal experiment.

TABLE 5. - $L = 4$, $i^* = 3$

Model	Param	After 500 observations				
1	θ_1	0	0	0	0	0
2	θ_2	0	0	0	0	0
3	θ_3	.953	.982	.953	.969	.954
4	θ_4	.047	.018	.047	.031	.046
1	β_0	0.1368	0.1385	0.1383	0.1394	0.1354
2	β_0	0.1376	0.1398	0.1391	0.1387	0.1369
	β_1	.2561	.2463	.2608	.2550	.2426
3	β_0	-0.0227	0.0085	-0.0069	-0.0028	-0.0096
	β_1	.2566	.2469	.2611	.2546	.2434
	β_2	.2906	.2392	.2648	.2548	.2677
4	β_0	-0.0227	0.0085	-0.0069	-0.0028	-0.0096
	β_1	.2355	.2385	.2399	.2715	.2225
	β_2	.2905	.2392	.2649	.2548	.2678
	β_3	.0281	.0112	.0281	-.0223	.0277
	p_0	0.178	0.178	0.178	0.178	0.178
	p_1	0	0	0	0	0
	p_2	.322	.320	.320	.318	.318
	p_3	0	0	0	.002	.004
	p_4	0	.006	.004	0	0
	p_5	.004	.004	.004	.002	.010
	p_6	.002	0	0	0	.004
	p_7	.318	.318	.318	.320	.312
	p_8	0	0	0	0	0
	p_9	.176	.174	.176	.180	.174

The values of the posterior probabilities and parameter means after 5 simulations, of 500 observations each, of the sequential selection procedure. Also listed are the proportions p_i of the times each $a(i)$ was chosen as the optimal experiment.

TABLE 6. - $L = 6$, $i^* = 3$

Model	Param	After 100 observations					After first 100 of 500 observations				
1	θ_1	0	0	0	0	0	0	0	0	0	0
2	θ_2	0	0	0	0	0	0	0	0	0	0
3	θ_3	.9554	.9564	.7630	.9098	.9432	.2733	.9534	.9471	.9449	.9554
4	θ_4	.0378	.0365	.1756	.0738	.0469	.5534	.0400	.0444	.0457	.0378
5	θ_5	.0052	.0055	.0451	.0130	.0075	.1168	.0050	.0070	.0075	.0052
6	θ_6	.0016	.0017	.0162	.0034	.0024	.0564	.0016	.0015	.0019	.0016
1	β_0	0.1567	0.1026	0.0891	0.0839	0.1259	*	*	*	*	*
2	β_0	0.1298	0.1197	0.1192	0.1118	0.1462	*	*	*	*	*
	β_1	.2696	.2258	.2151	.2422	.2527					
3	β_0	-0.0114	0.0098	-0.0332	0.0004	0.0156	*	*	*	*	*
	β_1	.2564	.2354	.2396	.2542	.2607					
	β_2	.2882	.2282	.3290	.2329	.2463					
4	β_0	-0.0110	0.0098	-0.0317	-0.0008	0.0161	*	*	*	*	*
	β_1	.2421	.2427	.3033	.2117	.2864					
	β_2	.2877	.2282	.3268	.2346	.2453					
	β_3	.0189	-.0099	-.0846	.0562	-.0341					
5	β_0	-0.0056	0.0200	-0.0082	-0.0175	0.0314	*	*	*	*	*
	β_1	.2458	.2403	.2870	.2232	.2855					
	β_2	.2562	.1692	.1788	.3427	.1641					
	β_3	.0148	-.0080	-.0676	.0432	.0336					
	β_4	.0271	.0512	.1288	-.0966	.0694					
6	β_0	-0.0061	0.0209	-0.0130	-0.0166	0.0300	*	*	*	*	*
	β_1	.2429	.1945	.3250	.2050	.2278					
	β_2	.2583	.1767	.1921	.3418	.1767					
	β_3	.0265	.1913	-.2368	.1214	.1958					
	β_4	.0254	.0432	.1203	-.0968	.0584					
	β_5	-.0089	-.1555	.1320	-.0609	-.1739					
	p_0	0.13	0.15	0.18	0.17	0.17	0.18	0.10	0.11	0.17	0.13
	p_1	.02	.08	.03	.06	.10	.01	.04	.11	.12	.02
	p_2	.17	.24	.32	.27	.22	.32	.12	.12	.20	.17
	p_3	.08	.03	0	0	.02	0	.05	.02	.02	.08
	p_4	.06	.03	.02	.07	.03	.01	.07	.14	.07	.06
	p_5	.03	.03	.07	.03	.05	0	.11	.02	.04	.03
	p_6	0	.12	.14	.14	.02	.06	.02	.06	.07	0
	p_7	.29	.17	.11	.09	.21	.26	.16	.15	.15	.29
	p_8	.05	.02	.01	.05	.04	0	.16	.12	.04	.05
	p_9	.17	.13	.12	.12	.14	.16	.17	.15	.12	.17

* Not recorded.

The values of the posterior probabilities and parameter means after 10 simulations, of 100 observations each, of the sequential selection procedure. The last 5 columns are data from the first 100 observations of the 500 observation simulations tabulated in table 7. The posterior means were not recorded for these 5 cases. Also listed are the proportions of the times each $a^{(i)}$ was chosen as the optimal experiment.

TABLE 7. - $L = 6$, $i^* = 3$

Model	Param	After 500 observations				
1	θ_1	0	0	0	0	0
2	θ_2	0	0	0	0	0
3	θ_3	.6046	.9812	.9722	.9746	.8526
4	θ_4	.3388	.0175	.0257	.0230	.1316
5	θ_5	.0425	.0009	.0018	.0021	.0125
6	θ_6	.0141	.0003	.0003	.0003	.0032
1	β_0	0.1321	0.1378	0.1325	0.1349	0.1413
2	β_0	0.1356	0.1278	0.1247	0.1168	0.1383
	β_1	.2434	.2616	.2541	.2549	.2581
3	β_0	-0.0026	0.0067	-0.0066	-0.0021	-0.0029
	β_1	.2446	.2556	.2486	.2466	.2571
	β_2	.2542	.2336	.2605	.2615	.2622
4	β_0	-0.0027	0.0067	-0.0067	-0.0017	-0.0030
	β_1	.2076	.2515	.2340	.2341	.2872
	β_2	.2544	.2335	.2608	.2609	.2623
	β_3	.0491	.0055	.0194	.0167	-.0399
5	β_0	-0.0113	0.0077	0.0111	0.0063	-0.0062
	β_1	.2081	.2518	.2332	.2394	.2873
	β_2	.2926	.2279	.2845	.2061	.2774
	β_3	.0486	.0052	.0203	.0101	-.0401
	β_4	-.0300	.0049	-.0200	.0496	-.0121
6	β_0	-0.0118	0.0047	-0.0100	0.0061	-0.0073
	β_1	.2164	.3128	.2432	.2386	.2682
	β_2	.2949	.2455	.2801	.2070	.2819
	β_3	.0130	-.2307	-.0225	.0134	.0405
	β_4	-.0317	-.0098	-.0168	.0490	-.0155
	β_5	.0270	.1773	.0328	-.0026	-.0618
	P_0	0.178	0.136	0.148	0.116	0.168
	P_1	.004	.064	.028	.070	.012
	P_2	.314	.206	.198	.086	.282
	P_3	.004	.036	.092	.112	.018
	P_4	.010	.014	.038	.110	.014
	P_5	0	.092	.014	.024	.016
	P_6	.022	.004	.012	.022	.002
	P_7	.296	.190	.280	.288	.304
	P_8	0	.104	.026	.014	.010
	P_9	.172	.154	.164	.158	.174

The values of the posterior probabilities and parameter means after five simulations, of 500 observations each, of the sequential selection procedure. Also listed are the proportions of the times each $a(i)$ was chosen as the optimal experiment.

TABLE 8. - $L = 6$, $i^* = 5$

Model	Param	After 100 observations					After first 100 of 500 observations				
		θ_1	θ_2	θ_3	θ_4	θ_5	θ_1	θ_2	θ_3	θ_4	θ_5
1	θ_1	0	0	0	0	0	0	0	0	0	0
2	θ_2	0	0	0	.021	.042	0	0	.011	.003	.026
3	θ_3	.945	.942	.956	.895	.848	.943	.877	.852	.904	.857
4	θ_4	.043	.043	.038	.063	.042	.048	.106	.089	.070	.076
5	θ_5	.007	.012	.005	.015	.033	.007	.013	.035	.018	.028
6	θ_6	.004	.003	.001	.006	.035	.002	.004	.013	.006	.013
1	β_0	0.1502	-0.0299	0.0231	-0.0079	0.0034	*	*	*	*	*
2	β_0	0.0356	0.0189	0.0288	0.0316	0.0431	*	*	*	*	*
	β_1	.5159	.5101	.5123	.5079	.5106					
3	β_0	-0.0348	-0.0413	-0.0412	-0.0098	0.0026	*	*	*	*	*
	β_1	.5040	.5077	.5096	.5084	.5133					
	β_2	.1467	.1265	.1478	.0837	.0850					
4	β_0	-0.0333	-0.0396	-0.0414	-0.0070	0.0026	*	*	*	*	*
	β_1	.4874	.5288	.5019	.5513	.5146					
	β_2	.1450	.1252	.1478	.0810	.0849					
	β_3	.0217	-.0261	.0099	-.0456	-.0016					
5	β_0	-0.0232	-0.0242	-0.0431	-0.0041	0.0161	*	*	*	*	*
	β_1	.4985	.5201	.5039	.5516	.4861					
	β_2	.0642	-.0063	.1616	.0184	-.1260					
	β_3	.0101	-.0185	.0074	-.0463	.0235					
	β_4	.0774	.1223	-.0131	.0605	-.2068					
6	β_0	-0.0298	-0.0234	-0.0435	-0.0067	0.0135	*	*	*	*	*
	β_1	.6005	.5398	.5233	.5143	.3504					
	β_2	.1187	-.0120	.1655	.0517	-.0226					
	β_3	-.3886	-.1048	-.0708	.1797	.6034					
	β_4	.0292	.1270	-.0167	.0301	.1071					
	β_5	.3037	.0673	.0599	-.1893	-.4482					
	p_0	0.10	0.16	0.14	0.24	0.18	0.18	0.14	0.17	0.18	0.22
	p_1	.05	.06	.07	.02	.13	.25	.04	.02	.01	.01
	p_2	.08	.20	.15	.03	.03	.06	.17	.04	.25	.02
	p_3	.04	.05	.07	.07	.02	.01	.06	.01	.03	0
	p_4	.08	.13	.08	.26	.18	.16	.06	.10	.10	.05
	p_5	.08	.08	.10	.12	.19	.07	.17	.25	.06	.41
	p_6	.08	.04	.06	.02	.03	.03	0	0	.14	.01
	p_7	.11	.10	.09	.03	.05	.09	.19	.19	.05	.02
	p_8	.22	.02	.10	0	.01	.06	.03	.01	.01	0
	p_9	.16	.16	.14	.21	.18	.09	.14	.21	.17	.26

* Not recorded.

The values of the posterior probabilities and parameter means after 10 simulations, of 100 observations each, of the sequential selection procedure. The last 5 columns are data from the first 100 observations of the 500 observation simulations tabulated in table 9. Also listed are the proportions of the times each $a^{(i)}$ was chosen as the optimal experiment.

TABLE 9. - $L = 6$, $i^* = 5$

Model	Param	After 500 observations				
1	θ_1	0	0	0	0	0
2	θ_2	0	0	0	0	0
3	θ_3	.974	.882	.899	.976	.976
4	θ_4	.020	.024	.029	.021	.022
5	θ_5	.003	.075	.062	.002	.002
6	θ_6	.002	.020	.009	0	0
1	β_0	-0.1051	0.0255	-0.0926	0.0130	0.0390
2	β_0	0.0290	0.0373	0.0436	0.0351	0.0458
	β_1	.4860	.5032	.4903	.5038	.5053
3	β_0	-0.0181	-0.0137	0.0018	-0.0258	0.0065
	β_1	.5016	.5008	.5019	.5056	.5061
	β_2	.1075	.1046	.0803	.1189	.0791
4	β_0	-0.0179	-0.0142	0.0027	-0.0257	0.0066
	β_1	.5035	.4859	.5201	.5177	.4956
	β_2	.1071	.1050	.0782	.1187	.0790
	β_3	-.0025	.0198	-.0237	-.0160	.0138
5	β_0	-0.0100	0.0012	0.0159	-0.0202	0.0136
	β_1	.4959	.4940	.5016	.5169	.4962
	β_2	.0413	-.0152	-.0328	.0897	.0306
	β_3	.0061	.0081	0	-.0151	.0129
	β_4	.0653	.1185	.1116	.0241	.0455
6	β_0	-0.0166	-0.0044	0.0157	-0.0186	0.0135
	β_1	.4232	.4352	.4617	.5039	.4949
	β_2	.0891	.0094	-.0250	.0835	.0309
	β_3	.2801	.2087	.1232	.0400	.0179
	β_4	.0224	.0996	.1030	.0288	.0452
	β_5	-.2070	-.1452	-.0854	-.0430	-.0039
	p_0	0.158	0.106	0.208	0.174	0.132
	p_1	.252	.190	.314	.010	.128
	p_2	.108	.054	.016	.306	.138
	p_3	.040	.052	.014	.006	.008
	p_4	.170	.184	.100	.020	.104
	p_5	.022	.034	.114	.030	.088
	p_6	.056	0	0	.072	.002
	p_7	.074	.134	.052	.208	.172
	p_8	.036	.118	.110	.022	.092
	p_9	.084	.128	.072	.152	.132

The values of the posterior probabilities and parameter means after five simulations, of 500 observations each, of the sequential selection procedure. Also listed are the proportions of the times each $a^{(i)}$ was chosen as the optimal experiment.

TABLE 10. - LARGE SAMPLE STUDY TWO

Model	Param	$\vec{\theta}_0 = (0.1, 0.2, 0.3, 0.4)$						$\vec{\theta}_0 = (0.25, 0.25, 0.25, 0.25)$						$\vec{\theta}_0 = (0.4, 0.3, 0.2, 0.1)$					
		5 simulations						5 simulations						5 simulations					
							Average												Average
1	θ_1	0	0	0	0	0	0	0	0	0	0	0	0	0	0	0	0	0	0
2	θ_2	0	0	0	0	0	0	0	0	0	0	0	0	0	0	0	0	0	0
3	θ_3	.957	.947	.828	.926	.960	.924	.969	.969	.966	.967	.930	.960	.984	.979	.976	.957	.984	.976
4	θ_4	.043	.053	.172	.074	.040	.076	.031	.031	.034	.033	.070	.040	.016	.021	.024	.043	.016	.024
1	θ_0	-0.0338	-0.0460	.0015	-0.0282	-0.0125	-0.0238	-0.0413	-0.0230	-0.0204	-0.0501	-0.0172	-0.0304	0.0392	0.0028	0.0045	-0.0195	0.0180	0.0090
2	θ_0	-0.0338	-0.0460	.0015	-0.0282	-0.0125	-0.0238	-0.0413	-0.0230	-0.0204	-0.0501	-0.0172	-0.0304	0.0392	0.0028	0.0045	-0.0195	0.0180	0.0090
	θ_1	1.0044	.9810	.9938	.9989	.9644	.9885	1.0142	.9779	1.0510	1.0487	.9697	1.012	.9484	1.0213	1.0603	.9906	.9614	.9964
3	θ_0	-0.0338	-0.0460	.0015	-0.0282	-0.0125	-0.0238	-0.0413	-0.0230	-0.0204	-0.0501	-0.0172	-0.0304	0.0392	0.0028	0.0045	-0.0195	0.0180	0.0090
	θ_1	1.0044	.9810	.9938	.9989	.9644	.9885	1.0142	.9779	1.0510	1.0487	.9697	1.012	.9484	1.0213	1.0603	.9906	.9614	.9964
	θ_2	-1.0555	-1.0012	-.9594	-1.0026	-.9948	-1.003	-.9753	-1.0205	-1.0208	-1.0139	-.9762	-1.001	-.9451	-.9539	-1.0253	-.9339	-1.0012	-.9719
4	θ_0	-0.0338	-0.0460	.0015	-0.0282	-0.0125	-0.0238	-0.0413	-0.0230	-0.0204	-0.0501	-0.0172	-0.0304	0.0392	0.0028	0.0045	-0.0195	0.0180	0.0090
	θ_1	1.0044	.9810	.9938	.9989	.9644	.9885	1.0142	.9779	1.0510	1.0487	.9697	1.012	.9484	1.0213	1.0603	.9906	.9614	.9964
	θ_2	-1.0555	-1.0012	-.9594	-1.0026	-.9948	-1.003	-.9753	-1.0205	-1.0208	-1.0139	-.9762	-1.001	-.9451	-.9539	-1.0253	-.9339	-1.0012	-.9719
	θ_3	.0110	-.0236	.0564	.0355	-.0016	.0155	-.0012	-.0052	-.0152	-.0121	.0418	.0016	.0019	-.0239	.0290	-.0456	0.0087	-.0060

The values of the posterior probabilities and parameter means after five simulations, of 500 observations each, of the sequential selection procedure with three different prior distributions on the models.

TABLE 11. - LARGE SAMPLE STUDY THREE

Model	Parameter	$\tau = 100.0$							$\tau = 1.0$							$\tau = 0.01$						
		5 simulations							5 simulations							5 simulations						
		Avg							Avg							Avg						
1	θ_1	1	1	1	1	1	1	1	1	1	1	1	1	1	1	0.9979	0.8155	0.9999	0.8951	0.9990	0.9415	
2	θ_2	0	0	0	0	0	0	0	0	0	0	0	0	0	0	.0006	.1843	0	.0012	.0004	.0373	
3	θ_3	0	0	0	0	0	0	0	0	0	0	0	0	0	0	.0015	.0002	.0001	.1037	.0006	.0212	
1	θ_0	0.953	1.027	1.094	0.979	1.106	1.032	1.111	1.091	1.046	1.298	0.932	1.096	1.415	0.454	1.221	0.916	0.906	0.932			
	θ_1	1.069	1.046	.972	1.067	.979	1.027	1.009	.890	.957	.800	.874	.906	1.093	1.274	1.318	.632	1.156	1.097			
	θ_2	.984	1.039	.997	.907	1.033	.992	.862	.910	1.210	.822	1.160	.993	1.163	.542	1.607	1.119	1.193	1.125			
2	θ_0	1.230	1.254	1.298	1.262	1.295	1.268	1.135	1.075	1.567	1.281	0.527	1.117	1.413	0.454	0.754	0.914	0.904	0.888			
	θ_1	1.559	1.564	1.469	1.520	1.494	1.521	1.030	.911	1.009	.816	.907	.935	1.122	1.275	1.319	.673	1.172	1.112			
	θ_2	.770	.746	.702	.738	.705	.732	.039	.019	.039	-.020	.219	.059	.074	.154	.031	.201	.018	.096			
3	θ_0	1.742	1.774	1.788	1.754	1.796	1.771	1.122	1.108	1.002	1.306	0.940	1.096	1.415	0.418	1.221	0.916	0.906	0.973			
	θ_1	1.515	1.559	1.480	1.438	1.520	1.502	.882	.928	1.246	.842	1.159	1.011	1.184	.613	1.607	1.121	1.198	1.145			
	θ_2	.258	.226	.212	.246	.204	.229	-.048	.013	-.041	-.048	.057	-.013	.056	.301	.005	.086	.012	.092			
	Number of trials until $\theta_{1,j} = 1.0$	3	3	3	3	3	3	49	47	39	59	52	49.2									

The posterior probabilities and posterior parameter means after five simulations of the sequential selection procedure for three different values of τ . For $\tau = 100.0$ and $\tau = 1.0$ the number of trials until $\theta_{1,j} = 1.0$ (within the accuracy of the computer) is also tabulated. For $\tau = 0.01$ the values are based upon 1000 observations.

TABLE 12. - SMALL SAMPLE STUDY ONE (H_3 TRUE)[$J_{MAX} = 8$]

θ_m	τ	$\vec{\mu}_{3,0}$	PCS	ASN	Starting value for random seq.			
0.70	0.5	(0, 0)	0.133	6.36	*			
.70	.5	(0.5, 0.5)	.458	7.15	041	574	501	221
.70	.5	(1.0, 1.0)	.544	6.82	261	404	147	531
.70	.5	(1.5, 1.5)	.446	5.89	251	233	175	021
.80	.5	(0, 0)	.173	7.50	265	603	111	061
.80	.5	(0.5, 0.5)	.468	7.78	237	616	233	015
.80	.5	(1.0, 1.0)	.531	7.52	066	231	644	355
.80	.5	(1.5, 1.5)	.460	7.24	124	715	646	251
.90	.5	(0, 0)	.229	7.98	202	255	025	241
.90	.5	(0.5, 0.5)	.479	7.92	020	625	757	465
.90	.5	(1.0, 1.0)	.513	7.81	154	510	176	555
.90	.5	(1.5, 1.5)	.439	7.76	043	355	261	141
.70	1.0	(0, 0)	.397	5.49	031	264	722	101
.70	1.0	(0.5, 0.5)	.673	5.88	142	153	215	611
.70	1.0	(1.0, 1.0)	.737	5.29	025	206	250	121
.70	1.0	(1.5, 1.5)	.621	4.84	244	233	735	061
.80	1.0	(0, 0)	.558	6.90	337	020	177	205
.80	1.0	(0.5, 0.5)	.755	6.94	361	341	044	651
.80	1.0	(1.0, 1.0)	.771	6.50	231	737	436	405
.80	1.0	(1.5, 1.5)	.700	6.22	107	152	460	271
.90	1.0	(0, 0)	.605	7.80	316	753	345	645
.90	1.0	(0.5, 0.5)	.765	7.45	042	264	053	551
.90	1.0	(1.0, 1.0)	.777	7.15	304	456	707	705
.90	1.0	(1.5, 1.5)	.689	7.12	324	670	521	455
.70	2.0	(0, 0)	.699	4.24	034	773	264	025
.70	2.0	(0.5, 0.5)	.871	4.03	361	656	711	721
.70	2.0	(1.0, 1.0)	.877	3.62	110	151	661	121
.70	2.0	(1.5, 1.5)	.723	3.48	000	766	306	641
.80	2.0	(0, 0)	.868	5.45	020	542	277	271
.80	2.0	(0.5, 0.5)	.962	4.99	073	755	766	635
.80	2.0	(1.0, 1.0)	.970	4.63	013	527	071	701
.80	2.0	(1.5, 1.5)	.872	4.61	041	554	522	311
.90	2.0	(0, 0)	.944	6.46	001	231	353	331
.90	2.0	(0.5, 0.5)	.967	5.66	361	503	245	415
.90	2.0	(1.0, 1.0)	.969	5.48	151	650	040	041
.90	2.0	(1.5, 1.5)	.939	5.80	233	434	565	701

*Not recorded.

Resulting PCS and ASN values for $J_{MAX} = 8$ and the combinations of θ_m , τ , and $\vec{\mu}_{3,0}$. Results are based upon 1500 simulations of the procedure for each combination.

TABLE 13. - SMALL SAMPLE STUDY ONE (H_3 TRUE)[$J_{MAX} = 16$]

θ_m	τ	$\mu_{3,0}$	PCS	ASN	Starting value for random seq.			
0.70	0.5	(0, 0)	0.354	9.48	272	036	225	461
.70	.5	(0.5, 0.5)	.665	10.7	057	343	345	741
.70	.5	(1.0, 1.0)	.723	9.63	073	144	502	151
.70	.5	(1.5, 1.5)	.555	7.38	231	500	657	525
.80	.5	(0, 0)	.508	13.6	033	254	034	051
.80	.5	(0.5, 0.5)	.761	13.3	225	553	740	341
.80	.5	(1.0, 1.0)	.806	12.3	134	537	257	651
.80	.5	(1.5, 1.5)	.661	11.8	251	356	646	745
.90	.5	(0, 0)	.574	15.5	056	537	424	615
.90	.5	(0.5, 0.5)	.752	14.6	246	632	674	651
.90	.5	(1.0, 1.0)	.800	13.9	140	077	157	311
.90	.5	(1.5, 1.5)	.710	13.8	044	035	362	005
.70	1.0	(0, 0)	.548	6.53	173	052	463	251
.70	1.0	(0.5, 0.5)	.821	6.82	063	364	104	441
.70	1.0	(1.0, 1.0)	.825	6.09	233	034	770	255
.70	1.0	(1.5, 1.5)	.637	5.36	017	237	125	325
.80	1.0	(0, 0)	.808	9.48	275	264	535	015
.80	1.0	(0.5, 0.5)	.971	9.30	015	352	360	531
.80	1.0	(1.0, 1.0)	.961	8.16	017	142	770	505
.80	1.0	(1.5, 1.5)	.865	7.86	004	724	275	765
.90	1.0	(0, 0)	.927	12.1	161	027	043	101
.90	1.0	(0.5, 0.5)	.973	10.8	101	732	737	651
.90	1.0	(1.0, 1.0)	.964	10.1	016	351	614	135
.90	1.0	(1.5, 1.5)	.958	10.6	171	716	572	235
.70	2.0	(0, 0)	.700	4.25	073	021	660	321
.70	2.0	(0.5, 0.5)	.878	4.17	003	466	340	375
.70	2.0	(1.0, 1.0)	.855	3.59	337	170	131	645
.70	2.0	(1.5, 1.5)	.714	3.51	055	666	256	215
.80	2.0	(0, 0)	.911	5.67	037	537	412	725
.80	2.0	(0.5, 0.5)	.990	5.12	111	525	350	761
.80	2.0	(1.0, 1.0)	.988	4.84	003	413	673	201
.80	2.0	(1.5, 1.5)	.894	4.71	055	643	644	455
.90	2.0	(0, 0)	.996	7.13	374	543	153	375
.90	2.0	(0.5, 0.5)	1.00	6.25	133	225	727	441
.90	2.0	(1.0, 1.0)	1.00	5.94	332	405	117	171
.90	2.0	(1.5, 1.5)	.995	6.20	010	312	536	461

Resulting PCS and ASN values for $J_{MAX} = 16$ and the combinations of θ_m , τ , and $\mu_{3,0}$. Results based upon 1000 simulations.

TABLE 14. - SMALL SAMPLE STUDY TWO (H_2 TRUE)[$J_{MAX} = 8$]

θ_m	τ	$\vec{\mu}_{2,0}$	PCS	ASN	Starting value for random seq.			
0.70	0.5	(1.0)	0.760	7.86	052	516	237	355
.80	.5	(1.0)	.734	7.98	016	160	602	721
.90	.5	(1.0)	.740	7.98	245	577	171	655
.70	1.0	(.5)	.828	7.63	321	722	414	631
.70	1.0	(1.0)	.882	7.20	340	321	470	071
.70	1.0	(1.5)	.800	6.86	360	415	546	645
.80	1.0	(1.0)	.872	7.98	273	760	237	431
.90	1.0	(.5)	.880	7.97	006	761	404	325
.90	1.0	(1.0)	.898	7.98	331	151	347	271
.90	1.0	(1.5)	.832	7.99	372	024	174	011
.70	2.0	(1.0)	.900	5.13	004	415	604	245
.80	2.0	(1.0)	.936	7.89	063	456	575	211
.90	2.0	(1.0)	.934	7.98	065	654	616	225

The PCS and ASN values resulting from 500 simulations of the sequential procedure for each of the tabulated combinations of θ_m , τ , and $\vec{\mu}_{2,0}$.

TABLE 15. - SMALL SAMPLE STUDY THREE (FOUR MODEL PROBLEM)

[J_{MAX} = 8]

θ_m	τ	$\vec{\mu}_{3,0}$	PCS	ASN	Starting value for random seq.			
0.70	1.0	(0, 1, -1)	0.767	7.55	*			
.70	1.0	(0, 0.5, 0)	.442	7.10	006	171	767	411
.70	1.0	(0, 0, 0.5)	.027	5.74	113	071	707	045
.70	1.0	(1, 0.5, 0)	.154	6.40	032	457	065	345
.80	1.0	(0, 1, -1)	.792	7.91	315	037	701	221
.80	1.0	(0, 0.5, 0)	.524	7.67	070	131	010	071
.80	1.0	(0, 0, 0.5)	.030	6.73	044	541	754	365
.80	1.0	(1, 0.5, 0)	.200	7.31	034	264	602	535
.90	1.0	(0, 1, -1)	.790	7.99	175	260	740	521
.90	1.0	(0, 0.5, 0)	.506	7.91	000	247	732	655
.90	1.0	(0, 0, 0.5)	.025	7.51	276	504	634	101
.90	1.0	(1, 0.5, 0)	.210	7.81	243	240	621	255

* Not recorded.

PCS and ASN values resulting from 1000 simulations performed for the indicated combinations of θ_m , τ , and $\vec{\mu}_{3,0}$.

TABLE 16. - CODED DATA FOR SAMPLE PROBLEM

(DATA TAKEN FROM DRAPER AND SMITH)

z_1	z_2	z_3	z_4	y
-75	0	0	-65	1.4
175	0	0	150	26.3
0	0	-65	150	29.4
0	0	165	-65	9.7
0	0	0	150	32.9
-75	-75	0	150	26.4
175	175	0	-65	8.4
-75	-75	-65	150	28.4
175	175	165	-65	11.5
0	0	-65	-65	1.3
0	0	165	150	21.4
0	-75	-65	-65	.4
0	175	165	150	22.9
0	0	0	-65	3.7
0	-75	0	150	26.5
0	-75	0	150	23.4
0	-75	0	150	26.5
0	175	0	-65	5.8
0	175	0	-65	7.4
0	175	0	-65	5.8
0	-75	-65	150	28.8
0	-75	-65	150	26.4
0	175	165	-65	11.8
0	175	165	-65	11.4

TABLE 17. - SUMMARY OF ANALYSIS OF EQUATION (7-1)

USING DATA OF TABLE 16

Term of model	Estimated coefficient	t-statistic	Descriptive significance level
z_4	0.112	28.4	0.999+
$z_3 z_4$	$-.354E-3$	6.8	.999+
z_2	$.323E-1$	3.4	.986
z_3	$.235E-1$	2.5	.955
z_1	$.319E-1$	2.1	.920
$z_1 z_4$	$-.416E-3$	1.9	.890
z_2^2	$.705E-4$	1.9	.886
z_3^2	$-.128E-3$	1.6	.836
$z_1 z_2$	$-.339E-3$	1.2	.717
$z_1 z_3$	$-.669E-4$.9	.576
z_1^2	$.705E-4$.5	.367
$z_2 z_3$	$.332E-4$.5	.347
$z_2 z_4$	$.178E-4$.3	.217
$R^2 = 0.988$ Residual mean square = 3.25 Replication mean square = 1.85 $F = \frac{\text{Lack-of-fit mean square}}{\text{Replication mean square}} = 2.90$			

TABLE 18a. - THE MODELS FOR EXAMPLE 1

$$H_1: y = \beta_0 + \beta_1 z_4 + \beta_2 z_3 z_4 + \epsilon$$

$$H_2: y = \beta_0 + \beta_1 z_4 + \beta_2 z_3 z_4 \\ + \beta_3 z_1 + \beta_4 z_2 + \beta_5 z_3 + \epsilon$$

$$H_3: y = \beta_0 + \beta_1 z_4 + \beta_2 z_3 z_4 \\ + \beta_3 z_1 + \beta_4 z_2 + \beta_5 z_3 \\ + \beta_6 z_2^2 + \beta_7 z_3^2 + \beta_8 z_1 z_4 + \epsilon$$

$$H_4: y = \beta_0 + \beta_1 z_4 + \beta_2 z_3 z_4 \\ + \beta_3 z_1 + \beta_4 z_2 + \beta_5 z_3 \\ + \beta_6 z_2^2 + \beta_7 z_3^2 + \beta_8 z_1 z_4 \\ + \beta_9 z_1^2 + \beta_{10} z_1 z_2 + \beta_{11} z_1 z_3 + \beta_{12} z_2 z_3 \\ + \beta_{13} z_2 z_4 + \beta_{14} z_4^2 + \epsilon$$

TABLE 18b. - THE PRIOR MEANS FOR EXAMPLE 1

$$\vec{\mu}_{1,0} = \begin{pmatrix} 1.215 \times 10 \\ 9.791 \times 10^{-2} \\ -2.650 \times 10^{-4} \end{pmatrix}$$

$$\vec{\mu}_{2,0} = \begin{pmatrix} 1.064 \ 10 \\ 1.113 \ 10^{-1} \\ -3.258 \ 10^{-4} \\ 2.211 \ 10^{-3} \\ 1.761 \ 10^{-2} \\ 1.066 \ 10^{-2} \end{pmatrix}$$

$$\vec{\mu}_{3,0} = \begin{pmatrix} 11.76 \\ .1137 \\ -.3376 \times 10^{-3} \\ .3322 \times 10^{-2} \\ .3114 \times 10^{-1} \\ .1768 \times 10^{-1} \\ -.1158 \times 10^{-3} \\ -.6788 \times 10^{-4} \\ -.1076 \times 10^{-3} \end{pmatrix}$$

$$\vec{\mu}_{4,0} = \begin{pmatrix} 12.70 \\ 0.1119 \\ -.3542 \times 10^{-3} \\ .3194 \times 10^{-1} \\ .3226 \times 10^{-1} \\ .2354 \times 10^{-1} \\ -.1512 \times 10^{-3} \\ -.1277 \times 10^{-3} \\ -.4164 \times 10^{-3} \\ .7045 \times 10^{-4} \\ -.3393 \times 10^{-3} \\ -.6690 \times 10^{-4} \\ .3323 \times 10^{-4} \\ .1785 \times 10^{-4} \\ 0 \end{pmatrix}$$

TABLE 18c. - THE MATRIX FROM WHICH $\Psi_{1,0}$ MAY BE TAKEN

Row 1	24.00000					
Row 2	1020.000	320700.0				
Row 3	-23950.00	0.381425E+07	0.210118E+10			
Row 4	300.0000	-14125.00	-0.114562E+07	108750.0		
Row 5	800.0000	-127250.0	577500.0	72500.00	290000.0	
Row 6	600.0000	-23950.00	0.307475E+07	33750.00	135000.0	188700.0
Row 7	290000.0	-0.380000E+07	-0.368156E+09	0.987500E+07	0.395000E+08	0.187500E+08
Row 8	188700.0	0.307475E+07	0.506112E+08	0.444750E+07	0.177900E+08	0.253050E+08
Row 9	-14125.00	0.172438E+07	0.231684E+09	0.193437E+07	-0.229375E+07	-0.114562E+07
Row 10	108750.0	0.193437E+07	-0.383297E+09	0.148125E+08	0.987500E+07	0.468750E+07
Row 11	72500.00	-0.229375E+07	-0.383297E+09	0.987500E+07	0.987500E+07	0.468750E+07
Row 12	33750.00	-0.114562E+07	-0.357216E+09	0.468750E+07	0.468750E+07	0.444750E+07
Row 13	13500.0	577500.0	-0.336394E+09	0.468750E+07	0.187500E+08	0.177900E+08
Row 14	-127250.0	-0.301625E+07	0.136534E+10	-0.229375E+07	-0.380000E+07	577500.0
Row 15	0	0	0	0	0	0
Row 7	0.775625E+10					
Row 8	0.343012E+10	0.455431E+10				
Row 9	-0.823281E+09	-0.357216E+09	0.122473E+10			
Row 10	0.193906E+10	0.857531E+09	0.804688E+07	0.290859E+10		
Row 11	0.193906E+10	0.857531E+09	-0.823281E+09	0.193906E+10	0.193906E+10	
Row 12	0.911719E+09	0.806719E+09	-0.383297E+09	0.911719E+09	0.911719E+09	0.857531E+09
Row 13	0.364687E+10	0.322687E+10	-0.383297E+09	0.911719E+09	0.911719E+09	0.857531E+09
Row 14	-0.205016E+10	-0.336394E+09	0.511906E+09	-0.823281E+09	-0.823281E+09	-0.383297E+09
Row 15	0	0	0	0	0	0
Row 13	0.343012E+10					
Row 14	-0.368156E+09	0.250450E+10				
Row 15	0	0	1.0			